Notes on Elementary Linear Analysis

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During the spring of 2021 I wrote an additional note to illustrate how the the spectral theorem for compact self-adjoint operators on Hilbert spaces may be applied on regular Sturm-Liouville problems. It is now included as the final chapter. The notes have otherwise been revised in accordance with a long list of suggestions I received from Ole Brevig, partly based on the feedback he got from some of the students enrolled in MAT3400/4400 (spring 2021). A warm thank goes to Ole and these students for their valuable contribution towards a better end result.

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CHAPTER 1

On normed spaces and bounded linear operators

1.1 Preliminaries

In this section we fix some notation and review some of the concepts and results that we will need. These are usually covered in undergraduate courses in real analysis, and the reader may consult the book of T. Lindstrøm, *Spaces: an introduction to real analysis* (AMS 2017), or any other standard book in real analysis, for examples, and proofs.

Throughout these notes \mathbb{F} will denote either \mathbb{R} (the real numbers) or \mathbb{C} (the complex numbers). If X, Y are sets, we let $X \times Y$ denote their Cartesian product, i.e.,

$$X \times Y = \left\{ (x, y) : x \in X, y \in Y \right\}.$$

A metric space (X, d) is called *complete* when every Cauchy sequence in (X, d) is convergent.

Definition 1.1.1. A normed space $(X, \|\cdot\|)$ over \mathbb{F} is a vector space X over \mathbb{F} which is equipped with a norm $\|\cdot\|$. We recall that X is then a metric space with respect to the metric given by $d(x, y) = \|x - y\|$ for $x, y \in X$. We will only consider normed spaces over \mathbb{F} in these notes, and we will often just write X to denote such a normed space, assuming tacitly that some norm on X is given.

When $x \in X$ and r > 0, we let $B_r^X(x)$ denote the closed ball in X with center in x and radius r, that is,

$$B_r^X(x) := \{ y \in X : \|x - y\| \le r \}.$$

When there is no danger of confusion, we just write $B_r(x)$ instead of $B_r^X(x)$. We also set

$$X_1 := B_1^X(0), \text{ i.e., } X_1 = \left\{ x \in X : ||x|| \le 1 \right\}.$$

Definition 1.1.2. If $(X, \|\cdot\|)$ is a normed space, and $\|\cdot\|'$ is also a norm on X, we say that $\|\cdot\|$ and $\|\cdot\|'$ are *equivalent* when there exist positive real numbers K and L such that

 $||x||' \leq K ||x||$ and $||x|| \leq L ||x||'$ for all $x \in X$.

When $\|\cdot\|$ and $\|\cdot\|'$ are equivalent, it is clear that a sequence $\{x_n\}_{n=1}^{\infty}$ in X converges to $x \in X$ w.r.t. $\|\cdot\|$ if and only if it converges to $x \in X$ w.r.t. $\|\cdot\|'$. The following proposition implies that for many purposes the choice of a norm in a finite-dimensional space can be made arbitrarily.

Proposition 1.1.3. If X is a finite-dimensional vector space, then all norms on X are equivalent.

Definition 1.1.4. Assume $\{x_n\}_{n=1}^{\infty}$ is a sequence in a normed space $(X, \|\cdot\|)$. We say that the series $\sum_{n=1}^{\infty} x_n$ is *convergent in* X if there is some $x \in X$ such that $\|x - \sum_{n=1}^{N} x_n\| \to 0$ as $N \to \infty$, in which case we say that $\sum_{n=1}^{\infty} x_n$ converges to x (w.r.t. $\|\cdot\|$), and also write $x = \sum_{n=1}^{\infty} x_n$.

Definition 1.1.5. When a normed space $(X, \|\cdot\|)$ is complete with respect to the associated metric given by

$$d(x,y) = \|x - y\|$$

for all $x, y \in X$, we say that X is a *Banach space*.

To check that a normed space is a Banach space, the following result is often useful:

Theorem 1.1.6. Let $(X, \|\cdot\|)$ be a normed space. Then X is a Banach space if and only if every absolutely convergent series in X is convergent in X, that is, if and only if the following condition holds :

Whenever $\sum_{n=1}^{\infty} x_n$ is a series in X such that $\sum_{n=1}^{\infty} ||x_n|| < \infty$, then $\sum_{n=1}^{\infty} x_n$ is convergent in X.

Remark 1.1.7. It is good to know that if X is a normed space, then we can always form its *completion*; this means that whenever needed, we can assume that X sits as a dense subspace of a Banach space \widetilde{X} where the norm of \widetilde{X} extends the norm on X. An elegant way to construct \widetilde{X} (as an application of the so-called Hahn-Banach theorem) is covered in more advanced courses on linear analysis.

Definition 1.1.8. Assume that X and Y are both vectors spaces over \mathbb{F} . Then a map $T: X \to Y$ is called a *linear operator* if we have

$$T(\lambda_1 x_1 + \lambda_2 x_2) = \lambda_1 T(x_1) + \lambda_2 T(x_2)$$

for all $\lambda_1, \lambda_2 \in \mathbb{F}$ and all $x_1, x_2 \in X$.

We let $\mathcal{L}(X, Y)$ denote the set of all linear operators from X to Y. One readily checks that $\mathcal{L}(X, Y)$ is a vector space over \mathbb{F} with respect to the operations defined by

$$(S+T)(x) = S(x) + T(x), \quad (\lambda T)(x) = \lambda T(x)$$

for $S, T \in \mathcal{L}(X, Y)$, $\lambda \in \mathbb{F}$ and $x \in X$. We also set $\mathcal{L}(X) := \mathcal{L}(X, X)$. We let $I_X \in \mathcal{L}(X)$ denote the *identity map from* X *into itself*, that is, $I_X(x) = x$ for all $x \in X$. We sometimes write I instead of I_X if no confusion is possible.

Definition 1.1.9. Assume that X and Y are both normed spaces over \mathbb{F} . Then a linear operator $T: X \to Y$ is called *bounded* if there exists some real number M > 0 such that

$$||T(x)|| \le M ||x|| \quad \forall x \in X,$$

or, equivalently, such that $||T(x)|| \leq M$ for all $x \in X_1$.

Proposition 1.1.10. Assume that X and Y are both normed spaces over \mathbb{F} and let $T \in \mathcal{L}(X, Y)$. Then the following conditions are equivalent:

- (a) T is bounded.
- (b) T is uniformly continuous on X.
- (c) T is continuous on X.
- (d) T is continuous at x = 0.

We will denote the set of all bounded linear operators from X to Y by $\mathcal{B}(X, Y)$. We follow tradition here and use the qualifying adjective "bounded", although we could equally well have used "continuous" instead. One readily checks that $\mathcal{B}(X, Y)$ is a subspace of $\mathcal{L}(X, Y)$. We also set $\mathcal{B}(X) = \mathcal{B}(X, X)$. **Proposition 1.1.11.** Assume that X and Y are both normed spaces over \mathbb{F} . For $T \in \mathcal{B}(X, Y)$, set

 $||T|| := \sup \{ ||T(x)|| : x \in X_1 \} < \infty.$

Then the map $T \mapsto ||T||$ is a norm on $\mathcal{B}(X, Y)$, called the operator norm, making $\mathcal{B}(X, Y)$ a normed space over \mathbb{F} . Moreover, we have

 $||T|| = \sup \{ ||T(x)|| : x \in X, ||x|| = 1 \} \quad (when \ X \neq \{0\}),$

and

$$||T(x)|| \le ||T|| ||x||$$
 for all $x \in X$.

Theorem 1.1.12. Assume that X is a normed space and Y is a Banach space (both over \mathbb{F}). Then $\mathcal{B}(X, Y)$ is Banach space. In particular, $\mathcal{B}(X)$ is a Banach space whenever X is a Banach space.

An immediate consequence of this theorem is that $\mathcal{B}(X, \mathbb{F})$ is a Banach space whenever X is normed space over \mathbb{F} . Elements of $\mathcal{L}(X, \mathbb{F})$ are called *linear functionals*. Thus $\mathcal{B}(X, \mathbb{F})$ consists of the *bounded linear functionals* on X; it is usually called the *dual space of* X and denoted by X^* in many books, or by X^{\sharp} in others.

Definition 1.1.13. A map $T: X \to Y$ between two vector spaces over \mathbb{F} is called a (*vector space*) *isomorphism* if $T \in \mathcal{L}(X, Y)$ and T is bijective (that is, T is both one-to-one and onto). It is then easy to check that the inverse map of $T, T^{-1}: Y \to X$, is linear, i.e., $T^{-1} \in \mathcal{L}(Y, X)$.

Definition 1.1.14. Assume that X and Y are normed spaces over \mathbb{F} . A map $T: X \to Y$ is called an *isomorphism of normed spaces* if T is a (vector space) isomorphism such that both T and T^{-1} are bounded.

Definition 1.1.15. Assume that X is a normed space and $T \in \mathcal{B}(X)$. Then we say that T is *invertible in* $\mathcal{B}(X)$ if T is an isomorphism of normed spaces. In other words, an operator $T \in \mathcal{B}(X)$ is invertible in $\mathcal{B}(X)$ if T is bijective and $T^{-1} \in \mathcal{B}(X)$.

Proposition 1.1.16. Let X, Y, Z be normed spaces over \mathbb{F} , and let $T \in \mathcal{B}(X,Y), S \in \mathcal{B}(Y,Z)$. Set $ST := S \circ T : X \to Z$. Then $ST \in \mathcal{B}(X,Z)$ and

$$||ST|| \leq ||S|| ||T||.$$

Corollary 1.1.17. Assume that X is a normed space and $S \in \mathcal{B}(X)$. For each $n \in \mathbb{N}$, let $S^n := S \cdots S$ denote the product of S with itself n times. Then $S^n \in \mathcal{B}(X)$ and

$$\|S^n\| \le \|S\|^n.$$

Note that by setting $S^0 = I_X$, this formula also holds when n = 0.

1.2 Norms and seminorms

Seminorms on vector spaces are almost as good as norms, and there is a standard way to produce a normed space from a vector space equipped with a seminorm. As we will use this procedure in the next chapter, we briefly describe it here. It is frequently used in linear analysis.

Definition 1.2.1. A seminorm on a vector space V (over \mathbb{F}) is a function $v \mapsto ||v||'$ from V into $[0, \infty)$ which is homogeonous and satisfies the triangle inequality, i.e., we have

$$\|\lambda v\|' = |\lambda| \|v\|'$$
 and $\|v + w\|' \le \|v\|' + \|w\|'$

for all $v, w \in V$ and $\lambda \in \mathbb{F}$.

Clearly, a seminorm $\|\cdot\|'$ is a norm if it also satisfies that

$$v \in V$$
 and $||v||' = 0 \implies v = 0$.

Example 1.2.2. Let X be a vector space and $(Y, \|\cdot\|)$ be a normed space (both over \mathbb{F}). Pick $T \in \mathcal{L}(X, Y)$. For each $x \in X$, set

$$||x||_T := ||T(x)||.$$

Then it is almost immediate that the map $x \mapsto ||x||_T$ is a seminorm on X. Since

$$||x||_T = 0 \Leftrightarrow ||T(x)|| = 0 \Leftrightarrow T(x) = 0 \Leftrightarrow x \in \ker(T),$$

we see that $\|\cdot\|_T$ gives a norm on X if and only if ker $(T) = \{0\}$, i.e., if and only if T is 1-1.

Let V be a vector space (over \mathbb{F}) and $\|\cdot\|'$ be a seminorm on V. Define a relation \sim on V by setting

$$v \sim w \Leftrightarrow ||v - w||' = 0, \quad v, w \in V.$$

It is an easy exercice to check that \sim is an equivalence relation (cf. Exercise 1.1). We denote the equivalence class of $v \in V$ by [v], that is,

$$[v] := \Big\{ w \in V : v \sim w \Big\},$$

and set $\widetilde{V} := \{ [v] : v \in V \}$. Moreover, for $v, w \in V$, and $\lambda \in \mathbb{F}$, we set

$$[v] + [w] := [v + w], \quad \lambda [v] := [\lambda v], \quad \| [v] \| := \|v\|'.$$

It is somewhat tedious, but straightforward, to check that these operations on \tilde{V} are well-defined, that \tilde{V} is a vector space (over \mathbb{F}) and that $\|\cdot\|$ is a norm on \tilde{V} . We leave it as an exercise to provide some of the necessary details.

1.3 Aspects of finite dimensionality

Unless otherwise specified, we always assume that the space \mathbb{F}^n , $n \in \mathbb{N}$, is equipped with the Euclidean norm $\|\cdot\|_2$ given by

$$||x||_2 = (|x_1|^2 + \dots + |x_n|^2)^{1/2}$$
 for $x = (x_1, \dots, x_n) \in \mathbb{F}^n$

and with the metric induced by this norm. As we recalled in Section 1.1, all norms on a finite-dimensional vector space are equivalent. The usual way to prove this is to consider first \mathbb{F}^n and show that any other norm on \mathbb{F}^n is equivalent to $\|\cdot\|_2$. A crucial fact in the proof is that a subset of \mathbb{F}^n is compact (w.r.t. the metric associated with $\|\cdot\|_2$) if and only if it is closed and bounded. (We recall that a subset K of a metric space is called *compact* if every sequence in K has a subsequence converging to a point in K.) It will be useful for us to know that this characterization of compactness, sometimes called the *Heine-Borel property*, holds in any finite-dimensional normed space. We will need the following lemma.

Lemma 1.3.1. Let X and Y be finite-dimensional normed spaces. Assume that X and Y are isomorphic as vector spaces and let $T \in \mathcal{L}(X, Y)$ be an isomorphism. Then T is an isomorphism of normed spaces.

Proof. We have to show that T and T^{-1} are bounded. We denote the respective norms on X and Y by the same symbol $\|\cdot\|$. For $x \in X$ we set

$$||x||_T := ||T(x)||.$$

Since T is 1-1, we get from Example 1.2.2 that the map $x \mapsto \|\cdot\|_T$ is a norm on X. Since X is finite-dimensional, $\|\cdot\|_T$ is equivalent to $\|\cdot\|$. In particular, this means that there exists some C > 0 such that

$$||T(x)|| = ||x||_T \le C ||x||$$
 for all $x \in X$,

which shows that T is bounded. Similarly, by considering the norm on Y given by $||y||_{T^{-1}} := ||T^{-1}(y)||$ for $y \in Y$, one deduces that T^{-1} is also bounded.

Proposition 1.3.2. Let X be a finite-dimensional normed space. Then a subset K of X is compact (w.r.t. the metric induced by the given norm) if and only if K is closed and bounded.

Proof. Since a compact subset of a metric space is always closed and bounded, we only have to show the reverse implication. So let $K \subseteq X$ be closed and

bounded. We must show that K is compact. If $X = \{0\}$, this is obviously true, so we may assume that $m := \dim(X) \ge 1$. Let then $T : X \to \mathbb{F}^m$ denote the coordinate map w.r.t. some basis for X. Lemma 1.3.1 gives that T is an isomorphism of normed spaces. Set $K' := T(K) \subseteq \mathbb{F}^m$. Then K' is bounded (since T is bounded). Moreover, K' is closed. Indeed, as $K' = (T^{-1})^{-1}(K)$, this follows from the continuity of T^{-1} . By the Heine-Borel property of \mathbb{F}^m , we can conclude that K' is compact. As $K = T^{-1}(K')$ and T^{-1} is continuous, this implies that K is compact, as desired.

Since the unit ball X_1 of a normed space is closed and bounded we get:

Corollary 1.3.3. The unit ball X_1 of a finite-dimensional normed space X is compact.

We recall that a vector space is said to be *infinite-dimensional* if it is not finite-dimensional. We note that if X is an infinite-dimensional normed space, then X_1 is *not* compact. (See Exercise 1.2.) In particular, this implies that an infinite-dimensional normed space never has the Heine-Borel property.

Another property which is automatically satisfied for a finite-dimensional normed space is completeness:

Proposition 1.3.4. Let X be a finite-dimensional normed space. Then X is a Banach space.

Proof. We may clearly assume that $X \neq \{0\}$. To show that X is complete, we let $\{x_n\}_{n\in\mathbb{N}}$ be a Cauchy sequence in X and have to prove that it is convergent. As in the proof of Proposition 1.3.2, we can pick an isomorphism of normed spaces $T: X \to \mathbb{F}^m$, where $m = \dim(X)$. For each $n \in \mathbb{N}$, set $y_n := T(x_n)$. Since $\|y_n - y_k\|_2 = \|T(x_n - x_k)\|_2 \leq \|T\| \|x_n - x_k\|$ for all $k, n \in \mathbb{N}$, we see that $\{y_n\}_{n\in\mathbb{N}}$ is a Cauchy sequence in \mathbb{F}^m . Since \mathbb{F}^m is complete, there exists $y \in \mathbb{F}^m$ such that $\|y_n - y\|_2 \to 0$ as $n \to \infty$. Set $x := T^{-1}(y) \in X$. Then we get

$$||x_n - x|| = ||T^{-1}(y_n - y)|| \le ||T^{-1}|| ||y_n - y||_2 \to 0 \text{ as } n \to \infty.$$

Thus, $\{x_n\}_{n \in \mathbb{N}}$ is convergent, as desired.

Corollary 1.3.5. Assume M is a finite-dimensional subspace of a normed space X. Then M is closed in X.

Proof. Assume $\{x_n\}_{n\in\mathbb{N}} \subseteq M$ converges to $x \in X$. We have to show that $x \in M$. As M is complete by Proposition 1.3.4, and $\{x_n\}_{n\in\mathbb{N}}$ is a Cauchy sequence in M, it follows that $\{x_n\}_{n\in\mathbb{N}}$ converges to some $y \in M$. Thus we get that $x = \lim_{n\to\infty} x_n = y \in M$.

1. On normed spaces and bounded linear operators

Finite dimensionality has also some impact on linear operators.

Example 1.3.6. Let $m, n \in \mathbb{N}$ and let $T \in \mathcal{L}(\mathbb{F}^n, \mathbb{F}^m)$. Then T is bounded. Indeed, let $A = [a_{i,j}]$ denote the standard matrix of T. Then we have $T(x) = (F_1(x), \ldots, F_n(x))$, where $F_i(x) := \sum_{j=1}^n a_{i,j} x_j$ for each $i = 1, \ldots, m$ and $x = (x_1, \ldots, x_n) \in \mathbb{F}^n$. Since each component F_i is clearly a continuous function from \mathbb{F}^n to \mathbb{F} , we get that T is continuous, and therefore bounded.

More generally, we have:

Proposition 1.3.7. Let X and Y be normed spaces and let $T \in \mathcal{L}(X, Y)$. Assume that X is finite-dimensional. Then T is bounded.

Proof. By replacing Y with T(X) if necessary, we may assume that Y is finite-dimensional. Moreover, we may also assume that both X and Y are different from $\{0\}$. Set $n = \dim(X), m = \dim(Y)$, and let $C : X \to \mathbb{F}^n$, $D : Y \to \mathbb{F}^m$ be isomomorphims, which are then necessarily isomorphisms of normed spaces by Lemma 1.3.1. The composition $T' := D \circ T \circ C^{-1}$ is then a linear map from \mathbb{F}^n to \mathbb{F}^m , hence it is bounded by the previous example. It follows that $T = D^{-1} \circ T' \circ C$, being the composition of bounded maps, is bounded.

Note that the above result is not true in general if we instead assume that Y is finite-dimensional, even in the case where $Y = \mathbb{F}$: a linear functional $T: X \to \mathbb{F}$ may be unbounded when X is an infinite-dimensional normed space. For an example, see Exercise 1.3.

Definition 1.3.8. A linear operator $T: X \to Y$ between two vector spaces X and Y is said to have *finite-rank* if the range of T is finite-dimensional, i.e., if dim $(T(X)) < \infty$.

It is obvious that a linear functional on a normed space has always finite-rank. As such a linear functional can be unbounded (cf. Exercise 1.3), this means that a finite-rank linear operator T between normed spaces is not necessarily bounded; in fact, it can be shown that a finite-rank linear operator T between normed spaces is bounded if and only if ker(T) is closed. Bounded finite-rank operators are examples of compact operators, which we will study in Chapter 4.

1.4 Extension by density and continuity

This short section is devoted to a very useful principle in linear analysis, often called *the principle of extension by density and continuity*. We will need the following elementary lemma, which is probably well-known.

Lemma 1.4.1. Assume that X and Y are metric spaces and f, g are continuous maps from X to Y which agree on a dense subset X_0 of X. Then f = g.

Proof. Let $x \in X$. Since X_0 is dense in X, there exists a sequence $\{x_n\}_{n \in \mathbb{N}}$ in X_0 which converges to x. By continuity of f and g, we get

$$f(x) = \lim_{n} f(x_n) = \lim_{n} g(x_n) = g(x).$$

Theorem 1.4.2. Assume that X is a normed space and Y is a Banach space (both over \mathbb{F}). Assume also that X_0 is a dense subspace of X, while Y_0 is a subspace of Y. Let $T_0 \in \mathcal{B}(X_0, Y_0)$. Then T_0 extends in a unique way to an operator $T \in \mathcal{B}(X, Y)$. It satisfies that $||T|| = ||T_0||$.

Proof. Let $x \in X$. Since X_0 is dense in X, there exists a sequence $\{x_n\}_{n\in\mathbb{N}}$ in X_0 such that $||x - x_n|| \to 0$ as $n \to \infty$. In particular, $\{x_n\}_{n\in\mathbb{N}}$ is a Cauchy sequence in X_0 . We claim that $\{T_0(x_n)\}_{n\in\mathbb{N}}$ is a Cauchy sequence in Y. Indeed, let $\varepsilon > 0$, and choose $N \in \mathbb{N}$ such that

$$||x_m - x_n|| < \varepsilon/||T_0||$$
 for all $m, n \ge N$.

Then, for all $m, n \in N$, we get

$$||T_0(x_m) - T_0(x_n)|| = ||T_0(x_m - x_n)|| = ||T_0|| ||x_m - x_n|| < \varepsilon,$$

as desired.

Since Y is complete, we can conclude that there exists some $y \in Y$ such that $\lim_n T_0(x_n) = y$. Note that y only depends on x. Indeed, assume $\{x'_n\}_{n\in\mathbb{N}}$ is another sequence in X_0 converging to x. Then the sequence $x_1, x'_1, x_2, x'_2, \ldots, x_n, x'_n, \ldots$ in X_0 also converges to x, so, arguing as above, we get that there exists some $z \in Y$ such that the sequence $T_0(x_1), T_0(x'_1), T_0(x_2), T_0(x'_2), \ldots, T_0(x_n), T_0(x'_n), \ldots$ converges to z. This implies that

$$\lim_{n} T_0(x'_n) = z = \lim_{n} T_0(x_n) = y.$$

Hence it makes sense to define T(x) := y. Doing this for every $x \in X$, we get a map $T : X \to Y$, and it is easy to check that T is linear, so we leave this as an exercise.

Next, we show that T is bounded. Let $x \in X$ and pick $\{x_n\}_{n \in \mathbb{N}}$ in X_0 converging to x. As $T(x) = \lim_n T_0(x_n)$ and $||T_0(x_n)|| \le ||T_0|| ||x_n||$ for all $n \in \mathbb{N}$, we get

$$||T(x)|| = \lim_{n} ||T_0(x_n)|| \le ||T_0|| \lim_{n} ||x_n|| = ||T_0|| ||x||.$$

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It follows that $T \in \mathcal{B}(X, Y)$ with $||T|| \leq ||T_0||$.

Further, T is an extension of T_0 . Indeed, let $x \in X_0$. Then set $x_n := x$ for all $n \in \mathbb{N}$. Since $\{x_n\}_{n \in \mathbb{N}}$ is a sequence in X_0 converging to x, we get that

$$T(x) = \lim_{x \to 0} T_0(x_n) = T_0(x).$$

The uniqueness of T as an extension of T_0 is immediate from Lemma 1.4.1. Finally, we have

$$||T_0|| = \sup\{||T_0(x)|| : x \in X_0, ||x|| \le 1\}$$

= sup{||T(x)|| : x \in X_0, ||x|| \le 1}
\le sup{||T(x)|| : x \in X, ||x|| \le 1} = ||T|| \le ||T_0||

Thus, $||T|| = ||T_0||$, as desired.

Remark 1.4.3. The conclusion of Theorem 1.4.2 is not necessarily true if Y is a normed space which is not complete (cf. Exercise 1.5).

An interesting special case of Theorem 1.4.2 is when T_0 is an isometry. We recall that a linear map between normed spaces is called an *isometry* when it is norm-preserving. A linear isometry is clearly bounded.

Corollary 1.4.4. Assume that X is a normed space, Y is a Banach space, X_0 is a dense subspace of X, Y_0 is a subspace of Y, and $U_0 \in \mathcal{L}(X_0, Y_0)$ is an isometry. Then the unique extension of U_0 to an operator U in $\mathcal{B}(X, Y)$ is also an isometry.

Proof. Theorem 1.4.2 guarantees that U_0 extends in a unique way to $U \in \mathcal{B}(X,Y)$. Let $x \in X$ and pick $\{x_n\}_{n \in \mathbb{N}}$ in X_0 converging to x. We then have $U(x) = \lim_n U_0(x_n)$, so we get

$$||U(x)|| = \lim_{n} ||U_0(x_n)|| = \lim_{n} ||x_n|| = ||x||.$$

Using Corollary 1.4.4, it can be shown that the completion of a (noncomplete) normed space is unique up to isometric isomorphism (cf. Exercise 1.7). We also record an important particular case of Theorem 1.4.2.

Corollary 1.4.5. Assume that X is a Banach space and X_0 is a dense subspace of X. Then every $T_0 \in \mathcal{B}(X_0)$ extends in a unique way to an operator $T \in \mathcal{B}(X)$, which satisfies that $||T|| = ||T_0||$.

An important application of this result to integral operators will be given in the next chapter (cf. Example 2.1.9).

1.5 Exercises

Exercise 1.1. Let V be a vector space (over \mathbb{F}) and let $\|\cdot\|'$ be a seminorm on V. In Section 1.2 we sketched how one can produce a normed space $(\tilde{V}, \|\cdot\|)$ from V and $\|\cdot\|'$. Here you are asked to provide some of the missing details.

a) Recall that \sim is defined by $v \sim w \Leftrightarrow ||v - w||' = 0$ $(v, w \in V)$. Check that \sim is an equivalence relation on your V, i.e., that the following properties hold for $u, v, w \in V$:

- $v \sim v$,
- $v \sim w \implies w \sim v$,
- $u \sim v$ and $v \sim w \implies u \sim w$.

b) For $v \in V$, set $[v] := \{ w \in V : v \sim w \}$. Further, for $v, w \in V$, and $\lambda \in \mathbb{F}$, set

 $[v] + [w] := [v + w], \quad \lambda [v] := [\lambda v], \quad || [v] || := ||v||'.$

Check that these operations on $\widetilde{V} := \{ [v] : v \in V \}$ are well-defined and turn \widetilde{V} into a vector space (over \mathbb{F}). Check also that $\| \cdot \|$ is a norm on \widetilde{V} .

Exercise 1.2. Let X be a normed space. Let M denote a finite-dimensional subspace of X, and assume $M \neq X$.

- a) Let $x \in X \setminus M$. Show that $d := \inf_{m \in M} ||x m|| > 0$.
- b) Show that there exists $y \in X$ such that ||y|| = 1 and

$$\frac{1}{2} \le \|y - m\| \quad \text{for all } m \in M.$$

c) Assume that X is infinite-dimensional (as a vector space). Show that the unit ball X_1 is not compact.

(*Hint*: Use b) to construct inductively a sequence $\{y_n\}_{n\in\mathbb{N}}$ in X_1 such that $1/2 \leq ||y_n - y_k||$ for all $1 \leq k < n$.)

Exercise 1.3. Recall that $\ell^{\infty}(\mathbb{N})$, which denotes the space of all bounded functions from \mathbb{N} into \mathbb{C} , is a normed space w.r.t. $||f||_u = \sup_{n \in \mathbb{N}} |f(n)|$. Let X be the subspace of $\ell^{\infty}(\mathbb{N})$ given by

 $X = \{ f : \mathbb{N} \to \mathbb{C} : f(n) = 0 \text{ for all but finitely many } n \}.$

a) Show that X is infinite-dimensional.

b) Consider X as a normed space w.r.t. $\|\cdot\|_u$ and let $L: X \to \mathbb{C}$ be defined by

$$L(f) = \sum_{n=1}^{\infty} f(n)$$

for all $f \in X$. Clearly, $L \in \mathcal{L}(X, \mathbb{C})$. Show that L is unbounded. Check also that ker(L) is not closed in X.

Exercise 1.4. Let $\mathcal{P}_{\mathbb{R}}$ denote the real vector space consisting of all polynomials in one real variable with real coefficients. For $p \in \mathcal{P}_{\mathbb{R}}$, set

$$||p|| := \sup_{t \in [0,1]} |p(t)|$$

- a) Explain why $p \to ||p||$ gives a well-defined norm on $\mathcal{P}_{\mathbb{R}}$.
- b) Define a linear operator $D: \mathcal{P}_{\mathbb{R}} \to \mathcal{P}_{\mathbb{R}}$ by

D(p) = p' (the derivative of p).

Show that D is unbounded. Conclude that $\mathcal{P}_{\mathbb{R}}$ is infinite-dimensional.

Exercise 1.5. Let X be a Banach space having a dense subspace X_0 which is not complete. Consider the identity map $I_0 : X_0 \to X_0$. Show that I_0 does not have an extension to a bounded linear map from X into X_0 .

Exercise 1.6. Assume that X is a normed space and Y is a Banach space (both over \mathbb{F}), and let $\{T_k\}_{k\in\mathbb{N}}$ be a sequence in $\mathcal{B}(X,Y)$ which is uniformly bounded in the sense that $M := \sup_{k\in\mathbb{N}} ||T_k|| < \infty$.

Moreover, assume that there exists a dense subset S of X such that $\{T_k(x)\}_{k\in\mathbb{N}}$ converges in Y for every $x \in S$.

Show that there exists $T \in \mathcal{B}(X, Y)$ such that

$$T(x) = \lim_{k} T_k(x)$$
 for all $x \in X$.

Exercise 1.7. Assume X_0 is a normed space and let (X, i) denote a *completion of* X_0 , that is, X is a Banach space and $i : X_0 \to X$ is a linear isometry such that $i(X_0)$ is dense in X. (As mentioned in Remark 1.1.7, such a completion always exists.)

Show that (X, i) is unique up to isometric isomorphism, meaning that the following holds: if (X', i') is another completion of X_0 , then there exists an isometric isomorphism $U: X \to X'$ such that $i' = U \circ i$.

CHAPTER 2

On L^p -spaces

An important class of Banach spaces over \mathbb{F} associated with measure spaces are the so-called L^p -spaces, where $1 \leq p \leq \infty$. We will assume that $\mathbb{F} = \mathbb{C}$, and just mention that the case where $\mathbb{F} = \mathbb{R}$ may be handled in a similar way. Our presentation is somewhat more detailed than the one given Lindstrøm's book *Spaces*. We assume that the reader is familiar with the basics of measure and integration theory, as covered for example in sections 7.1-7.6 and 8.1-8.4 of *Spaces*, supplied with Brevig's lecture note entitled *A measure* of Lebesgue measure. Our notation will essentially be the same as in these references.¹

2.1 The case $1 \le p < \infty$

Let (X, \mathcal{A}, μ) be a measure space and set

$$\mathcal{M} = \mathcal{M}(X, \mathcal{A}) := \left\{ f : X \to \mathbb{C} : f \text{ is } \mathcal{A}\text{-measurable} \right\}.$$

We recall that a complex function $f: X \to \mathbb{C}$ is called \mathcal{A} -measurable if it its real part and its imaginary part are \mathcal{A} -measurable. It is straightforward to show that \mathcal{M} is a vector space over \mathbb{C} (with its natural operations), so we leave this as an exercise for the reader (Exercise 2.1). We will be interested in subspaces of \mathcal{M} associated with any $p \in [1, \infty]$. In this section we consider the case $1 \leq p < \infty$.

¹In particular, if A is a subset of a set X, then $\mathbf{1}_A$ will denote the indicator function of A (in X). In most situations, the fact that X is not included in this notation will not create any confusion; a more precise notation could be $\mathbf{1}_A^X$. Indicator functions are sometimes called characteristic functions, and some authors prefer to write χ_A instead of $\mathbf{1}_A$. The σ -algebra on X consisting of all subsets of X will be denoted by $\mathcal{P}(X)$.

2. On L^p -spaces

Let $f \in \mathcal{M}$. Obviously, the function $|f|^p$ is non-negative, and one easily checks that it belongs to \mathcal{M} (using that the function $z \to |z|^p$ is continuous on \mathbb{C}). Thus, with the convention that $\infty^{1/p} := \infty$, we can define

$$||f||_p := \left(\int_X |f|^p \, d\mu\right)^{1/p} \in [0,\infty].$$

Moreover, we set

$$\mathcal{L}^p(X, \mathcal{A}, \mu) := \{ f \in \mathcal{M} : \|f\|_p < \infty \}.$$

We will just write \mathcal{L}^p , or $\mathcal{L}^p(\mu)$, or $\mathcal{L}^p(X)$, when it is clear from the context what we mean. We note that \mathcal{L}^1 consists of all the complex functions on X which are *integrable* (w.r.t. μ), i.e., which are \mathcal{A} -measurable and satisfies that $\int_X |f| d\mu < \infty$.

When $\mathcal{A} = \mathcal{P}(X)$ and μ is the counting measure on \mathcal{A} , it is common to write $\ell^p(X)$ instead of $\mathcal{L}^p(X, \mathcal{A}, \mu)$; the norm of $f \in \ell^p(X)$ is then given by $||f||^p = (\sum_{x \in X} |f(x)|^p)^{1/p}$. When $X = \mathbb{N}$, it is usual to write ℓ^p instead of $\ell^p(\mathbb{N})$, and think of elements of ℓ^p as sequences under the identification $f \to (f(1), f(2), f(3), \ldots)$.

It is not difficult to verify that \mathcal{L}^p is a subspace of \mathcal{M} . For example, closedness under addition follows readily from the inequality

$$|z+w|^p \le 2^p (|z|^p + |w|^p),$$

which is easily seen to hold for all $z, w \in \mathbb{C}$. On the other hand, it is not true in general that $\|\cdot\|_p$ is a norm on \mathcal{L}^p . The reason is that for $f \in \mathcal{L}^p$, we have

$$||f||_p = 0 \iff \int_X |f|^p \, d\mu = 0 \iff |f|^p = 0 \ \mu\text{-a.e.} \iff f = 0 \ \mu\text{-a.e.}$$

Using the triangle inequality for $|\cdot|$ on \mathbb{C} , one readily deduces that $||\cdot||_1$ gives a seminorm on \mathcal{L}^1 . As we will soon see, $||\cdot||_p$ is a seminorm on \mathcal{L}^p for any $p \geq 1$. To handle the case p > 1, we will need:

Theorem 2.1.1 (Hölder's inequality). Assume $p \in (1, \infty)$ and let $q \in (1, \infty)$ denote p's conjugate exponent, given by $q = \frac{p}{p-1}$, so that $\frac{1}{p} + \frac{1}{q} = 1$. Let $f \in \mathcal{L}^p$ and $g \in \mathcal{L}^q$. Then $fg \in \mathcal{L}^1$ and

$$||fg||_1 = \int_X |fg| \, d\mu \, \le \, ||f||_p \, ||g||_q \,. \tag{2.1.1}$$

Proof. We first note that if a, b are nonnegative real numbers, then we have

$$ab \leq \frac{a^p}{p} + \frac{b^q}{q}. \tag{2.1.2}$$

A geometric way to prove this inequality (called Young's inequality) is to observe that $A := \frac{a^p}{p} = \int_0^a x^{p-1} dx$ is the area of the region under the graph of the function $y = x^{p-1}$ over [0, a], while $B := \frac{b^q}{q} = \int_0^b y^{q-1} dy$ is the area of the region under the graph of the function $x = y^{q-1}$ over [0, b]. Now, as q - 1 = 1/(p-1), we have $y = x^{p-1} \Leftrightarrow x = y^{q-1}$ when $x, y \ge 0$. By making a drawing, one sees that ab, which is the area of the rectangle $[0, a] \times [0, b]$, is less than or equal to A + B, as desired.

Next, we note that we may assume that $||f||_p = ||g||_q = 1$. Indeed, assume that (2.1.1) holds whenever $||f||_p = ||g||_q = 1$, and consider $f \in \mathcal{L}^p$ and $g \in \mathcal{L}^q$. If $||f||_p = 0$ or $||g||_q = 0$, then both sides of (2.1.1) are equal to zero. On the other hand, if $||f||_p$ and $||g||_q$ are both nonzero, then we may use that (2.1.1) holds for the functions $f/||f||_p$ and $g/||g||_q$, and deduce that it holds in the general case.

Hence, assume that $||f||_p = ||g||_q = 1$. Then, using (2.1.2) with a = |f(x)|and b = |g(x)| for each $x \in X$, and linearity of the integral, we get

$$\begin{split} \int_{X} |fg| \, d\mu &= \int_{X} |f(x)| \, |g(x)| \, d\mu(x) \\ &\leq \frac{1}{p} \, \int_{X} |f(x)|^{p} \, d\mu(x) + \frac{1}{q} \, \int_{X} |g(x)|^{q} \, d\mu(x) \\ &= \frac{1}{p} \, \|f\|_{p}^{p} + \frac{1}{q} \, \|g\|_{q}^{q} \\ &= \frac{1}{p} + \frac{1}{q} = 1 \\ &= \|f\|_{p} \, \|g\|_{q} \,, \end{split}$$

as desired.

Corollary 2.1.2. Let $p \in [1, \infty)$. Then $\|\cdot\|_p$ is a seminorm on \mathcal{L}^p . In particular, for all $f, g \in \mathcal{L}^p$, we have

$$||f + g||_p \le ||f||_p + ||g||_p \quad (\text{Minkowski's inequality}) \tag{2.1.3}$$

Proof. As already mentioned, the case p = 1 is straightforward. So assume $p \in (1, \infty)$. The reader should have no problem to see that we have $\|\lambda f\|_p = |\lambda| \|f\|_p$ for all $\lambda \in \mathbb{C}$ and all $f \in \mathcal{L}^p$. Next, let $f, g \in \mathcal{L}^p$, and let q be p's conjugate exponent. As (p-1)q = p and p/q = p - 1, we have

$$\| |f + g|^{p-1} \|_{q} = \left(\int_{X} |f + g|^{(p-1)q} \, d\mu \right)^{1/q} = \left(\int_{X} |f + g|^{p} \, d\mu \right)^{1/q}$$
$$= \| f + g \|_{p}^{p/q} = \| f + g \|_{p}^{p-1} \, .$$

Since $f + g \in \mathcal{L}^p$, this shows that $|f + g|^{p-1} \in \mathcal{L}^q$; moreover, using Hölder's inequality (at the 4th step), we get

$$\begin{split} \|f+g\|_{p}^{p} &= \int_{X} |f+g|^{p} \, d\mu = \int_{X} |f+g| \, |f+g|^{p-1} \, d\mu \\ &\leq \int_{X} |f| \, |f+g|^{p-1} \, d\mu + \int_{X} |g| \, |f+g|^{p-1} \, d\mu \\ &\leq \|f\|_{p} \, \|\, |f+g|^{p-1}\|_{q} + \|g\|_{p} \, \|\, |f+g|^{p-1}\|_{q} \\ &= (\|f\|_{p} + \|g\|_{p}) \, \||f+g|^{p-1}\|_{q} \\ &= (\|f\|_{p} + \|g\|_{p}) \, \|f+g\|_{p}^{p-1}, \end{split}$$

and Minkowski's inequality clearly follows.

Let $\{f_n\}$ be a sequence in \mathcal{L}^p and $f \in \mathcal{L}^p$. We note that it may happen that $f_n \to f$ pointwise on X while $||f_n - f||_p \not\to 0$ as $n \to \infty$. For example one may let $X = \mathbb{R}$, $\mathcal{A} = \mathcal{B}_{\mathbb{R}}$, μ = Lebesgue measure on $\mathcal{B}_{\mathbb{R}}$, and consider the sequence given by $f_n = \mathbf{1}_{[n,n+1]}$ for each $n \in \mathbb{N}$: it converges pointwise to 0 on \mathbb{R} as $n \to \infty$, and satisfies $||f_n||_p = 1$ for all $n \in \mathbb{N}$.

The following \mathcal{L}^p -version of Lebesgue's Dominated Convergence Theorem gives conditions ensuring that a pointwise limit is also convergent w.r.t. $\|\cdot\|_p$.

Proposition 2.1.3. Let $p \in [1, \infty)$ and $\{f_n\}_{n \in \mathbb{N}} \subseteq \mathcal{L}^p$. Assume that there exist some $g \in \mathcal{L}^p$ such that $|f_n| \leq g$ μ -a.e. for all $n \in \mathbb{N}$, and some $f \in \mathcal{M}$ such that $f_n \to f$ pointwise μ -a.e. on X.

Then $f \in \mathcal{L}^p$ and $||f_n - f||_p \to 0$ as $n \to \infty$.

Proof. The assumptions imply that $|f_n|^p \leq g^p \mu$ -a.e. for all $n \in \mathbb{N}$ and that $|f_n|^p \to |f|^p$ pointwise μ -a.e. on X. It follows that we $|f|^p \leq g^p \mu$ -a.e., so

$$\int_X |f|^p \, d\mu \le \int_X g^p \, d\mu < \infty \,,$$

hence $f \in \mathcal{L}^p$. Further, we get

$$|f_n - f|^p \le (|f_n| + |f|)^p \le (2g)^p = 2^p g^p \quad \mu\text{-a.e.}$$

and $|f_n - f|^p \to 0$ pointwise μ -a.e. on X. Since $2^p g^p \in \mathcal{L}^1$, we can apply (the complex version of) Lebesgue's Dominated Convergence Theorem and get

$$\lim_{n \to \infty} \int_X |f_n - f|^p \, d\mu = \int_X 0 \, d\mu = 0$$

which gives that $||f_n - f||_p \to 0$ as $n \to \infty$, as desired.

We note that if $p \in [1, \infty)$ and $\{h_k\}_{k \in \mathbb{N}}$ is a sequence in \mathcal{L}^p such that $\lim_{k\to\infty} \|h_k - h\|_p = 0$ for some $h \in \mathcal{L}^p$, then it is not necessarily true that $h_k \to h$ pointwise μ -a.e. on X. However, the following weaker statement holds:

Proposition 2.1.4. Let $p \in [1, \infty)$. Assume that $\{h_k\}_{k \in \mathbb{N}}$ is a sequence in \mathcal{L}^p such that $\lim_{n\to\infty} \|h_k - h\|_p = 0$ for some $h \in \mathcal{L}^p$. Then there exists a subsequence $\{h_{k_n}\}_{n\in\mathbb{N}}$ which converges to h pointwise μ -a.e. on X.

We postpone the proof of this result to after the proof of Theorem 2.1.5.

Let $p \in [1, \infty)$. Since $\|\cdot\|_p$ is a seminorm on $\mathcal{L}^p = \mathcal{L}(X, \mathcal{A}, \mu)$ (by Corollary 2.1.2), we may use the procedure outlined in Section 1.2 to obtain a normed space $L^p = L^p(X, \mathcal{A}, \mu)$. We will denote the norm on L^p by the same symbol $\|\cdot\|_p$. Concretely, $(L^p, \|\cdot\|_p)$ can be described as follows. We first note that if $f, g \in L^p$, then

$$f \sim g \Leftrightarrow ||f - g||_p = 0 \Leftrightarrow f = g \quad \mu\text{-a.e.}$$

Hence, we have $[f] := \{g \in \mathcal{L}^p : f = g \ \mu\text{-a.e.}\}, \ L^p := \{[f] : f \in \mathcal{L}^p\}, \text{ and }$

 $[f] + [g] := [f + g], \quad \lambda [f] := [\lambda f], \quad \| [f] \|_p := \| f \|_p$

for $f, g \in \mathcal{L}^p$ and $\lambda \in \mathbb{C}$. Then we have:

Theorem 2.1.5. Let $p \in [1, \infty)$. Then $(L^p, \|\cdot\|_p)$ is a Banach space.

Proof. Let $\{[f_n]\}_{n\in\mathbb{N}}\subseteq L^p$ be such that $\sum_{n=1}^{\infty}\|[f_n]\|_p<\infty$, i.e., such that

$$S := \sum_{n=1}^{\infty} \|f_n\|_p < \infty.$$

According to Theorem 1.1.6 we have to show that the series $\sum_{n=1}^{\infty} [f_n]$ is convergent in L^p . It suffices to show that there exists some $F \in \mathcal{L}^p$ such that $\lim_{N\to\infty} \|\sum_{n=1}^N f_n - F\|_p = 0$, because this will give that

$$\lim_{N \to \infty} \|\sum_{n=1}^{N} [f_n] - [F]\|_p = \lim_{N \to \infty} \|\left[\sum_{n=1}^{N} f_n - F\right]\|_p = \lim_{N \to \infty} \|\sum_{n=1}^{N} f_n - F\|_p = 0,$$

thus showing that $\sum_{n=1}^{\infty} [f_n]$ converges to [F] in L^p .

For each $N \in \mathbb{N}$, set $g_N := \sum_{n=1}^N |f_n|$. Also, let $g: X \to [0, \infty]$ be given by

$$g(x) := \sum_{n=1}^{\infty} |f_n(x)|$$
 for all $x \in X$.

2. On L^p -spaces

Clearly, the sequence $\{g_N^p\}_{N\in\mathbb{N}}$ of \mathcal{A} -measurable nonnegative functions is nondecreasing, and it converges pointwise to the \mathcal{A} -measurable function g^p on X. Further, using Minkowski's inequality, we get

$$||g_N||_p \le \sum_{n=1}^N ||f_n||_p = \sum_{n=1}^N ||f_n||_p \le S$$

for all $N \in \mathbb{N}$. Hence, using the Monotone Convergence Theorem, we get

$$\int_X g^p d\mu = \lim_{N \to \infty} \int_X g_N^p d\mu = \lim_{N \to \infty} \|g_N\|_p^p \le S^p < \infty.$$

Since $g^p \ge 0$, it follows from Exercise 7.5.6 in Lindstrøm's book that g^p is finite μ -a.e., hence that g is finite μ -a.e. This means that the series $\sum_{n=1}^{\infty} f_n(x)$ is absolutely convergent for every x belonging to some $E \in \mathcal{A}$ such that $\mu(E^c) = 0$. We may therefore define $F \in \mathcal{M}$ by

$$F(x) = \begin{cases} \sum_{n=1}^{\infty} f_n(x) & \text{if } x \in E, \\ 0 & \text{if } x \in E^c. \end{cases}$$

Setting $F_N := \sum_{n=1}^N f_n$, we then have $|F_N| \le g_N \le g \in \mathcal{L}^p$ for every $N \in \mathbb{N}$, and $F_N \to F$ pointwise μ -a.e. on X as $N \to \infty$. Proposition 2.1.3 gives now that $F \in \mathcal{L}^p$ and $\lim_{N\to\infty} ||F_N - F||_p = 0$, as we wanted to show.

As a spin-off of the proof above, we can now prove Proposition 2.1.4:

Proof of Proposition 2.1.4. Let $\{h_k\}_{k\in\mathbb{N}}$ be a sequence in \mathcal{L}^p such that $\lim_{k\to\infty} \|h_k - h\|_p = 0$ for some $h \in \mathcal{L}^p$. Set $k_0 := 0$ and $h_0 := 0$. As $\{h_k\}_{k\in\mathbb{N}}$ is a Cauchy sequence in \mathcal{L}^p , it is not difficult to show (cf. Exercise 2.7) that we can pick a subsequence $\{h_{k_n}\}_{n\in\mathbb{N}}$ such that

$$\sum_{n=1}^{\infty} \|h_{k_n} - h_{k_{n-1}}\|_p < \infty.$$

Set $f_n := h_{k_n} - h_{k_{n-1}}$ for each $n \in \mathbb{N}$. Then $\{f_n\}_{n \in \mathbb{N}}$ is a sequence in \mathcal{L}^p such that $\sum_{n=1}^{\infty} \|f_n\|_p < \infty$, and the proof of Theorem 2.1.5 gives that the sequence $\{F_N\}_{N \in \mathbb{N}}$, given by $F_N := \sum_{n=1}^N f_n$, converges to a certain function $F \in \mathcal{L}^p$ pointwise μ -a.e. on X, such that $\lim_{N\to\infty} \|F_N - F\|_p = 0$. Now, for every $N \in \mathbb{N}$, we have

$$F_N = \sum_{n=1}^N (h_{k_n} - h_{k_{n-1}}) = h_{k_1} + (h_{k_2} - h_{k_1}) + \dots + (h_{k_N} - h_{k_{N-1}}) = h_{k_N}.$$

So we get that $\lim_{N\to\infty} \|h_{k_N} - F\|_p = 0$. Since $\lim_{n\to\infty} \|h_n - h\|_p = 0$, we also have $\lim_{N\to\infty} \|h_{k_N} - h\|_p = 0$. This gives that $\|F - h\|_p = 0$, hence that $F = h \mu$ -a.e. Thus we can conclude that $\{h_{k_N}\}_{N\in\mathbb{N}} = \{F_N\}_{N\in\mathbb{N}}$ converges to h pointwise μ -a.e. on X, as desired.

It is often good to have a dense subspace available. For our first result in this direction, we recall that a complex function on X is called *simple* if it takes a finite number of values. When p = 1, this result is essentially the same as Lemma 2.2.2 in Brevig's note. It is not difficult to adapt its proof to handle the case where p > 1. We leave this an exercise (cf. Exercise 2.8).

Proposition 2.1.6. Let $p \in [1, \infty)$. Denote by \mathcal{E} the subspace of \mathcal{M} consisting of all \mathcal{A} -measurable simple complex functions on X, and by \mathcal{E}^0 the subspace of \mathcal{E} spanned by the family $\{\mathbf{1}_A : A \in \mathcal{A}, \mu(A) < \infty\}$.

Then we have $\mathcal{E}^0 = \mathcal{E} \cap \mathcal{L}^p$. Moreover, $[\mathcal{E}^0] := \{[s] : s \in \mathcal{E}^0\}$ is a dense subspace of $L^p(X, \mathcal{A}, \mu)$ with respect to $\|\cdot\|_p$.

Our next result is the L^p -version of Littlewood's second principle for functions on bounded intervals.

Proposition 2.1.7. Let $p \in [1, \infty)$ and $a, b \in \mathbb{R}$, a < b. Let \mathcal{A} be the σ -algebra of all Lebesgue measurable subsets of [a, b] and μ be the Lebesgue measure on \mathcal{A} . Recall that C([a, b]) denotes the space of all continuous complex functions on [a, b].

Then $\{[g] : g \in C([a, b])\}$ is a dense subspace of $L^p([a, b])$ with respect to $\|\cdot\|_p$.

Note: since the linear map $g \mapsto [g]$ from C([a,b]) into $L^p([a,b])$ is 1-1, it is common to consider C([a,b]) as a dense subspace of $L^p([a,b])$ via this identification.

Proof. It is routine to verify that $\{[g] : g \in C([a, b])\}$ is a subspace of $L^p([a, b])$. To show that it is dense, one can argue in a quite similar way as in the proof of Lemma 2.2.3 in Brevig's note, by approximating simple functions with step functions and using Proposition 2.1.6.

Alternatively, one can exploit the regularity properties of the Lebesgue measure on \mathcal{A} to show that if $A \in \mathcal{A}$ and $\delta > 0$, then there exists some $k \in C([a, b])$ such that $\|\mathbf{1}_A - k\|_p < \delta$. Combining this fact with Proposition 2.1.6 readily gives the desired assertion. We leave the details to Exercise 2.9).

The L^p -version of Littlewood's second principle for functions on \mathbb{R} is analogous to Theorem 2.2.1 in Brevig's note (where p = 1).

Proposition 2.1.8. Let $p \in [1, \infty)$. Let \mathcal{A} be the σ -algebra of all Lebesgue measurable subsets of \mathbb{R} and μ be the Lebesgue measure on \mathcal{A} . Let $C_c(\mathbb{R})$

denote the space of all continuous complex functions on \mathbb{R} having compact support.²

Then $\{[g] : g \in C_c(\mathbb{R})\}$ is dense in $L^p(\mathbb{R})$ with respect to $\|\cdot\|_p$.

Proof. Minor modifications of the proof of Theorem 2.2.1 in Brevig's note (or of the previous proof) suffice. We leave this to the reader.

We end this section with an important example showing the usefulness of the principle of extension by density and continuity.

Example 2.1.9. Let $a, b \in \mathbb{R}$, a < b, and equip the space C([a, b]) of all continuous complex functions on [a, b] with the norm $||f||_2 = (\int_a^b |f(s)|^2 ds)^{1/2}$. Considering $[a, b] \times [a, b]$ as a metric space w.r.t. the Euclidean metric inherited from \mathbb{R}^2 , let $K : [a, b] \times [a, b] \to \mathbb{C}$ be a continuous function.

We can then associate to K an *integral operator* T_K on C([a, b]) as follows.

Let $f \in C([a, b])$. Since the function $t \mapsto K(s, t) f(t)$ is continuous on [a, b] for each $s \in [a, b]$, we may define a function $T_K(f) : [a, b] \to \mathbb{C}$ by

$$[T_K(f)](s) = \int_a^b K(s,t) f(t) dt \quad \text{for all } s \in [a,b].$$

We leave it as an exercise to verify, using basic knowledge from real analysis, that $T_K(f)$ is a continuous function on [a, b], satisfying

$$||T_K(f)||_2 \le \left(\int_a^b \int_a^b |K(s,t)|^2 \, ds \, dt\right)^{1/2} ||f||_2$$

As the map $f \mapsto T_K(f)$ is then clearly linear, it follows that T_K is a bounded linear operator from C([a, b]) into itself.

Let now $L^2([a, b])$ denote the L^2 -space associated with the measure space $([a, b], \mathcal{A}, \mu)$, where μ is the Lebesgue measure on the σ -algebra \mathcal{A} of all Lebesgue measurable subsets of [a, b].

As we may identify C([a, b]) with a dense subspace of $L^2([a, b])$ (cf. Proposition 2.1.7), we get from Corollary 1.4.5 that T_K has a unique extension to a bounded operator on $L^2([a, b])$, also denoted by T_K . The function K is usually called the *kernel* of the integral operator T_K . We will come back to such integral operators later.

More generally, one may define integral operators associated with kernels K which are \mathcal{L}^2 -functions on $[a, b] \times [a, b]$ (with respect to the so-called

²We recall that a function $g : \mathbb{R} \to \mathbb{C}$ is said to have *compact support* if there exists some M > 0 such that g = 0 outside [-M, M].

Lebesgue product measure), but this requires a thorough knowledge of integration theory on product spaces. We also note that we could have defined integral operators on $L^p([a, b])$ for any $p \in [1, \infty)$, in a similar way as we did above for p = 2.

2.2 The case $p = \infty$

We now consider the case $p = \infty$. Let $\mathcal{F}(X)$ denote the vector space consisting of all complex functions on X (with its natural operations). By an algebra of complex functions on X, we will mean a subspace of $\mathcal{F}(X)$ which is also closed under pointwise multiplication. For example, $\mathcal{M} = \mathcal{M}(X, \mathcal{A})$ is an algebra of complex functions on X. Another natural algebra is the one consisting of those functions in \mathcal{M} which are *bounded*. We will actually be interested in a slightly larger algebra.

Definition 2.2.1. A function $f \in \mathcal{M}$ is said to be *essentially bounded* (w.r.t. μ) if there exists some real number M > 0 such that

 $|f| \leq M$ μ -a.e.,

in which case we set $||f||_{\infty} := \inf \{M > 0 : |f| \leq M \mu$ -a.e. $\}$.

Example 2.2.2. a) Asume $g \in \mathcal{M}$ is bounded and set $||g||_u := \sup_{x \in X} |g(x)|$. Then g is essentially bounded (w.r.t. μ), and we have

$$\|g\|_{\infty} \le \|g\|_u \, .$$

Indeed, we have $\mu(\{x \in X : |g(x)| > ||g||_u\}) = \mu(\emptyset) = 0$. This gives that $|g| \leq ||g||_u \ \mu$ -a.e., and both assertions follow readily.

We note that it may happen that $||g||_{\infty} < ||g||_{u}$. For example, consider the Borel function g on $X = \mathbb{R}$ given by $g = \mathbf{1}_{\{0\}}$; letting μ be the Lebsgue measure on $\mathcal{B}_{\mathbb{R}}$, we get

$$||g||_{\infty} = 0 < 1 = ||g||_{u}.$$

b) Consider $X = [0, \infty)$, $\mathcal{A} =$ the Borel subsets of X and $\mu =$ the Lebesgue measure on \mathcal{A} . Let $f \in \mathcal{M}$ be given by

$$f(x) = e^{ix} + \sum_{n=1}^{\infty} n \mathbf{1}_{\{2n\pi\}}(x), \quad x \ge 0.$$

Then $f(2k\pi) = k+1$ for every $k \in \mathbb{N}$, so f is unbounded. On the other hand, f is essentially bounded (w.r.t. μ), with $||f||_{\infty} = 1$, since $\mu(|f|^{-1}((M,\infty)))$ is equal to 0 if $M \ge 1$ and to ∞ if 0 < M < 1.

The following useful observation may seem obvious, but it requires a proof.

Lemma 2.2.3. Let $f \in \mathcal{M}$ be essentially bounded (w.r.t. μ). Then we have

$$|f| \le ||f||_{\infty} \quad \mu\text{-a.e.}$$
 (2.2.1)

Proof. Set $B := \{x \in X : |f(x)| > ||f||_{\infty}\} \in \mathcal{A}$ and assume (for contradiction) that $\mu(B) > 0$. For each $n \in \mathbb{N}$, set

$$B_n := \left\{ x \in X : |f(x)| > ||f||_{\infty} + \frac{1}{n} \right\} \in \mathcal{A}.$$

Clearly, $B_n \subseteq B_{n+1}$ for every n, and $B = \bigcup_{n=1}^{\infty} B_n$, so we have

$$\lim_{n \to \infty} \mu(B_n) = \mu(B) > 0.$$

Hence there must exist at least one $N \in \mathbb{N}$ such that $\mu(B_N) > 0$. Now, by definition of $||f||_{\infty}$, we can find M > 0 such that $||f||_{\infty} \leq M < ||f||_{\infty} + \frac{1}{N}$ and $|f| \leq M \mu$ -a.e. But this implies that $|f| \leq ||f||_{\infty} + \frac{1}{N} \mu$ -a.e., i.e., $\mu(B_N) = 0$, and we have reached a contradiction.

Using Lemma 2.2.3, it is straightforward to verify that the set $\mathcal{L}^{\infty} = \mathcal{L}^{\infty}(X, \mathcal{A}, \mu)$ consisting of all functions in \mathcal{M} that are essentially bounded (w.r.t. μ) is an algebra of complex functions on X (cf. Exercise 2.12). Another application is the following Hölder-type inequality:

Proposition 2.2.4. Let $q \in [1, \infty)$, $f \in \mathcal{L}^{\infty}$ and $g \in \mathcal{L}^{q}$. Then $fg \in \mathcal{L}^{q}$ and

$$||fg||_q \leq ||f||_{\infty} ||g||_q$$

Proof. Using Lemma 2.2.3 we get that $|fg|^q = |f|^q |g|^q \le ||f||_{\infty}^q |g|^q \mu$ -a.e. It follows that

$$\int_X |fg|^q \, d\mu \, \leq \, \|f\|_\infty^q \int_X |g|^q \, d\mu \, < \infty \, .$$

Hence $fg \in \mathcal{L}^q$. Moreover, taking the q-th root, we obtain the desired inequality.

Convergence in \mathcal{L}^{∞} with respect to $\|\cdot\|_{\infty}$ is closely related to uniform convergence:

Proposition 2.2.5. Let $\{f_n\}_{n\in\mathbb{N}} \subseteq \mathcal{L}^{\infty}$ and $f \in \mathcal{L}^{\infty}$. Then we have that $\|f_n - f\|_{\infty} \to 0$ as $n \to \infty$ if and only if there exists some $E \in \mathcal{A}$ such that $\mu(E^c) = 0$ and $f_n \to f$ uniformly on E.

Proof. Assume $||f_n - f||_{\infty} \to 0$ as $n \to \infty$. For each $n \in \mathbb{N}$, set

$$F_n := \{ x \in X : |f_n(x) - f(x)| > ||f_n - f||_{\infty} \} \in \mathcal{A}.$$

Since $f_n - f \in \mathcal{L}^{\infty}$, we have $\mu(F_n) = 0$ for each *n*. Hence, $F := \bigcup_{n \in \mathbb{N}} F_n \in \mathcal{A}$ and $\mu(F) = 0$.

Set now $E := F^c \in \mathcal{A}$. Then $\mu(E^c) = 0$ and

$$E = \{x \in X : |f_n(x) - f(x)| \le ||f_n - f||_{\infty} \text{ for all } n \in \mathbb{N}\}.$$

It is then obvious that $f_n \to f$ uniformly on E. The proof of the reverse implication goes along the same lines, and we leave it as an exercise (cf. Exercise 2.14).

As with the \mathcal{L}^p -spaces for $1 \leq p < \infty$, an annoying fact is that in general $\|\cdot\|_{\infty}$ is only a seminorm on \mathcal{L}^{∞} . But we can use again the procedure outlined in Section 1.2 to obtain a norm $\|\cdot\|_{\infty}$ on a vector space $L^{\infty} = L^{\infty}(X, \mathcal{A}, \mu)$, given as follows. Setting

$$[f] = \{g \in \mathcal{L}^{\infty} : \|f - g\|_{\infty} = 0\} = \{g \in \mathcal{L}^{\infty} : f = g \ \mu\text{-a.e.}\}$$

for each $f \in \mathcal{L}^{\infty}$, we get

$$L^{\infty} = L^{\infty}(X, \mathcal{A}, \mu) := \{ [f] : f \in \mathcal{L}^{\infty} \}$$

with operations $[f] + [g] := [f + g], \lambda[f] := [\lambda f]$, and norm $||[f]||_{\infty} := ||f||_{\infty}$ (where $f, g \in \mathcal{L}^{\infty}$ and $\lambda \in \mathbb{C}$).

As to be expected, we have:

Theorem 2.2.6. $(L^{\infty}, \|\cdot\|_{\infty})$ is a Banach space.

Proof. We have to show that L^{∞} is complete w.r.t. the metric associated with $\|\cdot\|_{\infty}$.

Let $\{[f_n]\}_{n\in\mathbb{N}}$ be a Cauchy sequence in L^{∞} . So each f_n belongs to \mathcal{L}^{∞} and for any given $\varepsilon > 0$, there exists some $N \in \mathbb{N}$ such that

$$m, n \ge N \implies || [f_m] - [f_n] ||_{\infty} < \varepsilon,$$

that is,

$$m, n \ge N \implies || f_m - f_n ||_{\infty} < \varepsilon.$$
 (2.2.2)

For each $m, n \in \mathbb{N}$, set

$$F_{m,n} := \left\{ x \in X : |f_m(x) - f_n(x)| > ||f_m - f_n||_{\infty} \right\}.$$

Then $F_{m,n} \in \mathcal{A}$ and $\mu(F_{m,n}) = 0$ for all $m, n \in \mathbb{N}$ (because $f_m - f_n \in \mathcal{L}^{\infty}$).

Next, set $F := \bigcup_{m,n \in \mathbb{N}} F_{m,n} \in \mathcal{A}$ and $E := F^c \in \mathcal{A}$.

Note that $\mu(E^c) = \mu(F) = 0$ (since $0 \le \mu(F) \le \sum_{m,n \in \mathbb{N}} \mu(F_{m,n}) = 0$). Moreover,

$$E = \bigcap_{m,n\in\mathbb{N}} (F_{m,n})^c = \bigcap_{m,n\in\mathbb{N}} \left\{ x \in X : |f_m(x) - f_n(x)| \le ||f_m - f_n||_{\infty} \right\}$$
$$= \left\{ x \in X : |f_m(x) - f_n(x)| \le ||f_m - f_n||_{\infty} \text{ for all } m, n \in \mathbb{N} \right\}.$$

Let now $\varepsilon > 0$ be given, and choose $N \in \mathbb{N}$ such that (2.2.2) holds.

Then for all $x \in E$ and all $m, n \geq N$, we have

$$|f_m(x) - f_n(x)| \le ||f_n - f_m||_{\infty} < \varepsilon.$$
 (2.2.3)

It follows that $\{f_n(x)\}_{n\in\mathbb{N}}$ is a Cauchy sequence in \mathbb{C} for each $x \in E$. Since \mathbb{C} is complete, this implies that $\{f_n(x)\}_{n\in\mathbb{N}}$ is convergent for each $x \in E$, hence that $\lim_{n\to\infty} f_n(x) = g(x)$ for some $g(x) \in \mathbb{C}$ for each $x \in E$. Thereby we obtain a function $g : E \to \mathbb{C}$, which is \mathcal{A}_E -measurable since g is the pointwise limit of the restriction of the f_n 's to E. (Here, \mathcal{A}_E denotes the σ -algebra of all sets in \mathcal{A} which are contained in E).

We can now extend g to an \mathcal{A} -measurable function $f: X \to \mathbb{C}$ by setting f(x) = g(x) if $x \in E$, and f(x) = 0 otherwise.

Again, let $\varepsilon > 0$ be given and choose N as above. Then, for all $x \in E$ and all $m \in \mathbb{N}$ such that $m \ge N$, we get from (2.2.3) that

$$|f_m(x) - f(x)| = |f_m(x) - g(x)| = \lim_{n \to \infty} |f_m(x) - f_n(x)| \le \varepsilon.$$

This implies that $\{f_m\}_{m \in \mathbb{N}}$ converges uniformly to f on E.

Moreover, set $D := E \cap \{x \in X : |f_N(x)| \leq \|f_N\|_{\infty}\} \in \mathcal{A}$. Then we have

$$|f(x)| = |f(x) - f_N(x) + f_N(x)| \le |f(x) - f_N(x)| + |f_N(x)| \le \varepsilon + ||f_N||_{\infty}$$

for all $x \in D$. As

$$0 \le \mu(D^c) \le \mu(F) + \mu(\{x \in X : |f_N(x)| > ||f_N||_{\infty}\}) = 0,$$

we have $\mu(D^c) = 0$, so

$$|f| \leq \varepsilon + ||f_N||_{\infty} \mu$$
-a.e.

This shows that $f \in \mathcal{L}^{\infty}$. Using Proposition 2.2.5, we can now conclude that $||f_m - f||_{\infty} \to 0$ as $m \to \infty$. Thus

$$\|[f_m] - [f]\|_{\infty} = \|f_m - f\|_{\infty} \to 0 \text{ as } m \to \infty.$$

This means that $\{[f_m]\}_{j\in\mathbb{N}}$ converges to [f] in L^{∞} . We have thereby shown that every Cauchy sequence in L^{∞} is convergent and the proof is finished.

2.3 Exercises

In the exercises of this section, unless otherwise specified, (X, \mathcal{A}, μ) denotes a measure space and \mathcal{M} denotes the space of \mathcal{A} -measurable complex functions on X.

Exercise 2.1. a) Show that \mathcal{M} is a vector space (over \mathbb{C}). You can take it as granted that the set $\mathcal{F}(X)$ consisting of all complex functions on X is a vector space (over \mathbb{C}) w.r.t. its natural pointwise defined operations.

b) Let $f \in \mathcal{M}$. Show that $|f| \in \mathcal{M}$, and that $|f|^p \in \mathcal{M}$ for every p > 1.

Exercise 2.2. Let $f \in \mathcal{L}^1$. Check that $\operatorname{Re}(f)$ and $\operatorname{Im}(f)$ are both integrable, and set

$$\int_X f \, d\mu := \int_X \operatorname{Re}(f) \, d\mu + i \, \int_X \operatorname{Im}(f) \, d\mu \, .$$

Then show that $\left| \int_X f \, d\mu \right| \leq \int_X |f| \, d\mu$.

Exercise 2.3. Assume that $X = [1, \infty)$, $\mathcal{A} =$ the Borel subsets of X and μ is the Lebesgue measure on \mathcal{A} . Let $f \in \mathcal{M}$ be given by

$$f(x) = \frac{1}{x}$$
 for all $x \ge 1$,

and let $1 \leq p < \infty$. Show that $f \in \mathcal{L}^p(X, \mathcal{A}, \mu)$ if and only if p > 1, and compute $||f||_p$ in this case.

Exercise 2.4. Assume that $X = \mathbb{R}$, $\mathcal{A} =$ the Borel subsets of X and μ is the Lebesgue measure on \mathcal{A} . Let $f \in \mathcal{M}$ be given by

$$f(x) = e^{-x^2}$$
 for all $x \in \mathbb{R}$.

Show that $f \in \mathcal{L}^p(X, \mathcal{A}, \mu)$ for all $p \in [1, \infty)$ and compute $||f||_p$. (You are allowed to use that $\lim_{N\to\infty} \int_{-N}^N e^{-t^2} dt = \sqrt{\pi}$ without proof.)

Exercise 2.5. Assume that X = (0, 1], $\mathcal{A} =$ the Borel subsets of X and μ is the Lebesgue measure on \mathcal{A} . Let $1 \leq p < \infty$ and $f \in \mathcal{M}$ be given by

$$f(x) = \frac{1}{\sqrt{x}}$$
 for all $x \in (0, 1]$.

a) Show that $f \in \mathcal{L}^p(X, \mathcal{A}, \mu)$ if and only if p < 2, and compute $||f||_p$ in this case.

b) Let ν be the measure on \mathcal{A} given by

$$\nu(A) = \int_A x \ d\mu(x) \quad \text{ for all } A \in \mathcal{A}.$$

Show that $f \in \mathcal{L}^p(X, \mathcal{A}, \nu)$ if and only if p < 4, and compute $||f||_p$ in this case.

Exercise 2.6. Assume that $X = [1, \infty)$, $\mathcal{A} =$ the Borel subsets of X and μ is the Lebesgue measure on \mathcal{A} . For each $n \in \mathbb{N}$, define $f_n \in \mathcal{M}$ by

$$f_n(x) = \frac{n}{n x^{1/3} + 1}$$
 for all $x \ge 1$.

a) Show that $f_n \in \mathcal{L}^p$ for all $n \in \mathbb{N}$ whenever 3 .

b) Assume that $3 . Decide whether the sequence <math>\{[f_n]\}_{n \in \mathbb{N}}$ is convergent in L^p and find its limit if it converges.

Exercise 2.7. Let $\{x_k\}_{k\in\mathbb{N}}$ be a sequence in a vector space X having a seminorm $\|\cdot\|$, and assume this sequence is Cauchy, i.e., for every $\varepsilon > 0$, there exists some $N_{\varepsilon} \in \mathbb{N}$ such that $\|x_k - x_l\| < \varepsilon$ for all $k, l \ge N_{\varepsilon}$. Set $k_0 := 0$ and $x_0 := 0$. Show that we can pick $k_1 < k_2 < \cdots < k_n < \cdots$ in \mathbb{N} such that

$$\sum_{n=1}^{\infty} \|x_{k_n} - x_{k_{n-1}}\| < \infty.$$

Exercise 2.8. Prove Proposition 2.1.6.

Exercise 2.9. Provide the details missing in the proofs of Proposition 2.1.7 and Proposition 2.1.8.

Exercise 2.10. Provide the details missing in Example 2.1.9.

Exercise 2.11. The purpose of this exercise is to present an alternative way to produce integral operators on $L^2([a, b])$ associated to continuous kernels.

Let μ denote the Lebesgue measure on the Lebesgue measurable subsets of [a, b] and let K be a continuous complex function on $[a, b] \times [a, b]$. For each $s \in [a, b]$, let $k_s : [a, b] \to \mathbb{C}$ denote the continuous function defined by

$$k_s(t) := K(s,t)$$
 for all $t \in [a,b]$.

a) Let $f \in \mathcal{L}^2([a, b])$ and $s \in [a, b]$. Show the function $k_s f$ is Lebesgue integrable on [a, b] and satisfies

$$\int_{[a,b]} k_s f \, d\mu \, \leq \|k_s\|_2 \, \|f\|_2 \, .$$

b) Let $f \in \mathcal{L}^2([a, b])$ and define $g : [a, b] \to \mathbb{K}$ by

$$g(s) = \int_{[a,b]} k_s f \, d\mu = \int_{[a,b]} K(s,t) f(t) \, d\mu(t) \quad \text{for each } s \in [a,b] \, .$$

Show that g is continuous and check that

$$||g||_2 \le M ||f||_2$$
, where $M := \left(\int_a^b \int_a^b |K(s,t)|^2 \, ds dt\right)^{1/2}$

Deduce that we obtain a linear map $T_K^0: \mathcal{L}^2([a,b]) \to \mathcal{L}^2([a,b])$ by setting

$$(T_K^0(f))(s) := \int_{[a,b]} k_s f d\mu$$
 for each $f \in \mathcal{L}^2([a,b])$ and all $s \in [a,b]$,

which satisfies that $||T_K^0(f)||_2 \leq M ||f||_2$ for all $f \in \mathcal{L}^2([a, b])$.

c) Check that the operator $T_K : L^2([a, b]) \to L^2([a, b])$ defined by

$$T_K([f]) = [T_K^0(f)] \text{ for all } [f] \in L^2([a, b])$$

is well-defined, linear and bounded, with $||T_K|| \leq M$.

Exercise 2.12. Check that $\|\cdot\|_{\infty}$ is a seminorm on \mathcal{L}^{∞} (so that $\|\cdot\|_{\infty}$ gives a norm on L^{∞}). Check also that \mathcal{L}^{∞} is an algebra of functions on X and that we have $\|fg\|_{\infty} \leq \|f\|_{\infty} \|g\|_{\infty}$ for all $f, g \in \mathcal{L}^{\infty}$.

Exercise 2.13. Let $f \in \mathcal{M}$. Show that $f \in \mathcal{L}^{\infty}$ if and only if there exists a bounded function $g \in \mathcal{M}$ such that $f = g \mu$ -a.e., in which case we have

$$||f||_{\infty} = \inf\{ ||g||_{u} : g \in \mathcal{M} \text{ is bounded and } g = f \ \mu\text{-a.e.} \}.$$

Exercise 2.14. Finish the proof of Proposition 2.2.5.

Exercise 2.15. Let $1 \le p \le r < \infty$ and X be a nonempty set. Show that

$$\ell^p(X) \subseteq \ell^r(X) \subseteq \ell^\infty(X).$$

Exercise 2.16. Let $p \in [1, \infty)$ and assume that the measure space (X, \mathcal{A}, μ) is *finite*, that is, $\mu(X) < \infty$.

- a) Show that $\mathcal{L}^{\infty} \subseteq \mathcal{L}^{p}$.
- b) Consider $1 \leq p \leq r < \infty$ and let $f \in \mathcal{L}^r$. Show that $f \in \mathcal{L}^p$ and

$$||f||_p \leq \mu(X)^{\frac{1}{p} - \frac{1}{r}} ||f||_r.$$

Hint: Use Hölder's inequality in a suitable way.

Note that this shows that $\mathcal{L}^r \subseteq \mathcal{L}^p$. In particular, we have $\mathcal{L}^\infty \subseteq \mathcal{L}^2 \subseteq \mathcal{L}^1$.

c) Consider the Lebesgue measure on the Borel subsets of \mathbb{R} . Give an example of a function which is in \mathcal{L}^2 , but not in \mathcal{L}^1 Give also an example of a function which is in \mathcal{L}^{∞} , but not in \mathcal{L}^2 .

Exercise 2.17. Let \mathcal{E} denote the space of all simple functions belonging to \mathcal{M} and let $f \in \mathcal{L}^{\infty}$.

Show that there exists a sequence $\{h_n\}_{n\in\mathbb{N}}$ in \mathcal{E} such that $||f - h_n||_{\infty} \to 0$ as $n \to \infty$. Deduce that the space $[\mathcal{E}] := \{[h] : h \in \mathcal{E}\}$ is dense in L^{∞} with respect to $|| \cdot ||_{\infty}$.

CHAPTER 3

On Hilbert spaces and bounded linear operators

3.1 Inner product spaces

We assume that inner product spaces over \mathbb{R} are familiar to the reader. In applications, one frequently has to work with inner products spaces over \mathbb{C} . We give here a unified review of the basics results about such spaces (but skip the proofs as they are almost identical in the complex case.)

Definition 3.1.1. An *inner product space over* \mathbb{F} is a vector space X over \mathbb{F} which is equipped with an inner product $\langle \cdot, \cdot \rangle : X \times X \to \mathbb{F}$. This means that for $x, y, z \in X$ and $\lambda \in \mathbb{F}$ we have:

- i) $\langle x + y, z \rangle = \langle x, z \rangle + \langle y, z \rangle$,
- ii) $\langle \lambda x, y \rangle = \lambda \langle x, y \rangle$,
- iii) $\langle y, x \rangle = \overline{\langle x, y \rangle},$
- iv) $\langle x, x \rangle \ge 0$,
- v) $\langle x, x \rangle = 0$ if and only if x = 0.

Remark 3.1.2. a) Properties i) and ii) say that the inner product is linear in the first variable.

b) When $\mathbb{F} = \mathbb{R}$, property iii) says that the inner product is symmetric, i.e., $\langle y, x \rangle = \langle x, y \rangle$; combining i) and ii) with iii), we then get that the inner product is also linear in the second variable.

c) When $\mathbb{F} = \mathbb{C}$, we get that the inner product is *conjugate-linear* in the second variable; this means that we have

$$\langle x, y + z \rangle = \langle x, y \rangle + \langle x, z \rangle$$
 and $\langle x, \lambda y \rangle = \lambda \langle x, y \rangle$.

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Some authors prefer to use inner products that are linear in the second variable and conjugate-linear in the first variable. This is common in textbooks related to physics or mathematical physics. As one can go from one type to the other by setting $\langle x, y \rangle' := \langle y, x \rangle$, it is mainly a matter of taste which convention one chooses to use.

Example 3.1.3. The standard example of an inner product space over \mathbb{F} is \mathbb{F}^n , where $n \in \mathbb{N}$, equipped with

$$\langle \mathbf{x}, \mathbf{y} \rangle := \sum_{j=1}^{n} x_j \overline{y_j}, \quad \mathbf{x} = (x_1, \cdots, x_n), \mathbf{y} = (y_1, \cdots, y_n) \in \mathbb{F}^n.$$

In the sequel, by an inner product space, we will always mean an inner product space over \mathbb{F} . An inequality of fundamental importance is:

Theorem 3.1.4 (The Cauchy-Schwarz inequality). Let X be an inner product space. Then we have

$$|\langle x, y \rangle| \le ||x|| \, ||y|| \tag{3.1.1}$$

for all $x, y \in X$, with equality if and only if x and y are linearly dependent.

If X is an inner product space, then using the Cauchy-Schwarz inequality, one deduces that $||x|| := \langle x, x \rangle^{1/2}$ gives a norm on X. Thus, X is then a normed space, and its norm is easily seen to satisfy the *parallellogram law*, that is, for all $x, y \in X$ we have

$$||x + y||^{2} + ||x - y||^{2} = 2 ||x||^{2} + 2 ||y||^{2}.$$
 (3.1.2)

Moreover, the so-called *polarization identities* are sometimes useful:

$$\langle x, y \rangle = \frac{1}{4} \left(\|x + y\|^2 - \|x - y\|^2 \right) \text{ if } \mathbb{F} = \mathbb{R},$$
 (3.1.3)

$$\langle x, y \rangle = \frac{1}{4} \sum_{k=0}^{3} i^k \|x + i^k y\|^2 \text{ if } \mathbb{F} = \mathbb{C}.$$
 (3.1.4)

Definition 3.1.5. Let X be an inner product space. If $x, y \in X$, then x and y are said to be *orthogonal* (to each other) when $\langle x, y \rangle = 0$. If U, V are nonempty subsets of X, we write $U \perp V$ when $\langle u, v \rangle = 0$ for all $u \in U, v \in V$. A nonempty subset S of X is called *orthogonal* if x and y are orthogonal for all $x, y \in S$ such that $x \neq y$. Moreover, S is called *orthonormal* if S is orthogonal and ||x|| = 1 for all $x \in S$.
Proposition 3.1.6 (*Pythagoras*). Assume $\{x_1, \ldots, x_n\}$ is a finite orthogonal subset of an inner product space X. Then we have

$$||x_1 + \dots + x_n||^2 = ||x_1||^2 + \dots + ||x_n||^2.$$

Proposition 3.1.7. Assume $S = \{u_1, \ldots, u_n\}$ is a finite orthonormal subset of an inner product space X. Then S is linearly independent. Moreover, if $u \in Span\{u_1, \ldots, u_n\}$, i.e., if u is a linear combination of the vectors in S, then

$$u = \sum_{j=1}^{n} \langle u, u_j \rangle u_j$$
 and $||u||^2 = \sum_{j=1}^{n} |\langle u, u_j \rangle|^2.$

Proposition 3.1.8 (Bessel's inequality). Assume $S = \{u_j : j \in J\}$ is an orthonormal subset of an inner product space X indexed by a countable set J.¹ Then for any $x \in X$ we have

$$\sum_{j \in J} |\langle x, u_j \rangle|^2 \le ||x||^2.$$

Definition 3.1.9. An inner product space X (over \mathbb{F}) is called an *Hilbert* space (over \mathbb{F}) when X is complete as a normed space (with respect to the norm associated with its inner product).

Example 3.1.10. Since a finite-dimensional normed space is automatically complete. we get that \mathbb{F}^n is a Hilbert space w.r.t. its standard inner product.

Remark 3.1.11. Assume X is an inner product space. Considering X as a normed space, we may form its completion \widetilde{X} (cf. Remark 1.1.7), and extend the inner product on X to an inner product on \widetilde{X} as follows: if $y, y' \in \widetilde{X}$, then we can pick sequences $\{x_n\}_{n=1}^{\infty}, \{x'_n\}_{n=1}^{\infty}$ in X converging respectively to y and y'; after checking that $\{\langle x_n, x'_n \rangle\}_{n=1}^{\infty}$ is a Cauchy sequence in \mathbb{F} , hence is convergent, we may set

$$\langle y, y' \rangle := \lim_{n \to \infty} \langle x_n, x'_n \rangle.$$

It is then a somewhat tedious exercise to verify that this gives a well-defined inner product on \widetilde{X} which extends the one on X. This means that whenever needed, we may assume that X sits as a dense subspace of a Hilbert space \widetilde{X} (called the *completion* of X) where the inner product on \widetilde{X} extends the inner product on X.

¹We recall that a set J is *countable* if it is either finite or countably infinite (that is, there exists a bijection from \mathbb{N} onto J).

We end this section with a rich class of examples.

Example 3.1.12. Let (X, \mathcal{A}, μ) be a measure space. Set $\mathcal{L}^2 := \mathcal{L}^2(X, \mathcal{A}, \mu)$ and $L^2 := L^2(X, \mathcal{A}, \mu)$. We can organize L^2 as a Hilbert space (over \mathbb{C}) as follows.

Let $f, g \in \mathcal{L}^2$. Then \overline{g} is measurable (since $\overline{g} = \operatorname{Re}(g) - i \operatorname{Im}(g)$) and $\int_X |\overline{g}|^2 d\mu = \int_X |g|^2 d\mu = ||g||_2^2 < \infty$, so $\overline{g} \in \mathcal{L}^2$. Hence, Hölder's inequality gives that $f \overline{g} \in \mathcal{L}^1$, and we can set

$$\langle [f], [g] \rangle := \int_X f \,\overline{g} \, \mathrm{d}\mu \,.$$

We leave it as an exercise to check that this gives a well-defined inner product on L^2 . As the associated norm obviously coincides with the $\|\cdot\|_2$ -norm, L^2 is complete w.r.t. this norm and we can conclude that L^2 is a Hilbert space.

3.2 Geometry in Hilbert spaces

In courses in elementary linear algebra, one learns that if M is a finitedimensional subspace of an inner product space H, then every vector in H can be written in a unique way as the sum of a vector in M and a vector in the orthogonal complement M^{\perp} . As we are going to establish, such a decomposition also holds when H is a Hilbert space and M is closed subspace of H, not necessarily finite-dimensional.

We recall first that if (X, d) is a metric space, $x \in X$ and A is a nonempty subset of X, then the *distance from* x to A is defined by

$$d(x, A) = \inf\{d(x, y) : y \in A\}.$$

If for example A is compact, then the function $y \mapsto d(x, y)$, being continuous, will attain its minimum on A; hence, in this case, there exists some (not necessarily unique) $x_A \in A$ such that $d(x, A) = d(x, x_A)$. However, if A is closed, but not compact, such an element x_A may not exist (cf. Exercise 3.3).

Let us now consider a Hilbert space H with the metric d_H associated to its norm. If $x \in H$ and M is a closed subspace of H, then M is not compact, so we can not use the result mentioned above. However, if M is finite-dimensional, then we know from previous courses that there exists a unique $x_M \in M$ which gives the best approximation to x in M, in the sense that

$$||x - x_M|| \le ||x - y|| \quad \text{for all } y \in M,$$

which means that $d_H(x, x_M) = d_H(x, M)$. Moreover, we also know that the vector x_M is constructed as the orthogonal projection of x on M. When M is closed, but not finite-dimensional, we are going to switch this process by showing first that there exists a unique best approximation x_M to x in M, and use this fact to define the orthogonal projection of x on M.

We will actually prove a more general result, valid for any closed convex subset of H. We recall that a subset C of some vector space V (over \mathbb{F}) is called *convex* if C contains the line segment between any two elements of C, i.e., if we have $(1 - t)x + ty \in C$ whenever $x, y \in C$ and $t \in [0, 1]$.

Clearly, any subspace of a vector space is convex, as is any ball in a normed space. Using that the norm in a Hilbert space satisfies the parallellogram law, we will prove the following result, which the reader is advised to illustrate geometrically by looking at various examples in \mathbb{R}^2 .

Theorem 3.2.1. Let C be a nonempty closed convex subset of a Hilbert space H and let $x \in H$. Then there is a unique vector $x_C \in C$ such that $d_H(x, x_C) = d_H(x, C)$, that is, such that

$$||x - x_C|| \le ||x - y|| \quad for \ all \ y \in C.$$

The vector x_C is called the best approximation to x in C.

Proof. We first consider the case where x = 0. We then have to show that there is a unique vector $0_C \in C$ of minimal norm, i.e., which satisfies that

$$||0_C|| = \inf \{ ||y|| : y \in C \}$$

Set $s := \inf \{ \|y\|^2 : y \in C \}$. For each $n \in \mathbb{N}$ we can find $y_n \in C$ such that

$$s \leq ||y_n||^2 < s + \frac{1}{2n}$$
 (3.2.1)

Then the sequence $\{y_n\}_{n \in \mathbb{N}}$ is Cauchy in *H*. Indeed, consider $m, n \in \mathbb{N}$. Then, using the parallellogram law and (3.2.1), we get that

$$||y_n + y_m||^2 + ||y_n - y_m||^2 = 2 ||y_n||^2 + 2 ||y_m||^2 < 4s + \frac{1}{n} + \frac{1}{m}$$

Now, since C is convex, we have $c := \frac{1}{2}y_n + \frac{1}{2}y_m \in C$. Hence,

$$||y_n + y_m||^2 = 4 ||c||^2 \ge 4s$$
,

so we get

$$||y_n - y_m||^2 < 4s + \frac{1}{n} + \frac{1}{m} - ||y_n + y_m||^2 \le \frac{1}{n} + \frac{1}{m}.$$

Thus, given $\varepsilon > 0$, we can choose $N \in \mathbb{N}$ such that $N \ge 2\varepsilon^{-2}$, and obtain that $||y_n - y_m|| < \varepsilon$ for all $n, m \ge N$, as desired.

As *H* is complete, there exists $y_0 \in H$ such that $\lim_n y_n = y_0$. Since *C* is closed, $y_0 \in C$. Letting $n \to \infty$ in (3.2.1), we get that

$$||y_0|| = \sqrt{s} = \inf\{||y|| : y \in C\}.$$

If $y'_0 \in C$ also satisfies that $||y'_0|| = \inf\{||y|| : y \in C\}$, then we can consider the sequence $\{z_n\}_{n\in\mathbb{N}}$ in C given by $z_n = y'_0$ if n is odd and $z_n = y_0$ if n is even. Since z_n satisfies (3.2.1) (with $y_n = z_n$) for each n, we can conclude as above that $\{z_n\}_{n\in\mathbb{N}}$ is convergent. This clearly implies that $y'_0 = y_0$.

Thus, y_0 is the unique vector in C satisfying $||y_0|| = \inf\{||y|| : y \in C\}$, and we can set $0_C := y_0$.

In the general case where x is any vector in H, we note that the set

$$D := \{x - y : y \in C\}$$

is closed and convex. Using the first part, we get that there exists a unique vector $0_D \in D$ such that $||0_D|| = \inf\{||z|| : z \in D\} = d(x, C)$. Then $x_C := x - 0_D \in C$ has the desired properties.

Since a subspace is convex, Theorem 3.2.1 gives:

Corollary 3.2.2. Let M be a closed subspace of a Hilbert space H and let $x \in H$. Then there is a unique vector $x_M \in M$, called the best approximation to x in M, satisfying that

$$||x - x_M|| \le ||x - y|| \quad for \ all \ y \in M.$$

This result has some far-reaching consequences. To help us formulate these, we first introduce some appropriate terminology.

Definition 3.2.3. Let M_1 and M_2 be subspaces of a Hilbert space H. We will say that H is the *algebraic direct sum of* M_1 and M_2 if every $x \in H$ can be written in a unique way as $x = m_1 + m_2$ with $m_1 \in M_1$ and $m_2 \in M_2$.

Moreover, we will say that H is the *direct sum of* M_1 and M_2 , and write

$$H = M_1 \oplus M_2 \,,$$

when M_1 and M_2 are both closed, and H is the algebraic direct sum of M_1 and M_2 .

Our main interest in this course is the concept of direct sum (which, by the way, clearly also makes sense in a normed space). However, the following characterization of algebraic direct sum will be useful. **Lemma 3.2.4.** Let M_1 and M_2 be subspaces of a Hilbert space H and set

$$M_1 + M_2 := \{ m_1 + m_2 : m_1 \in M_1, \, m_2 \in M_2 \}.$$

Then the following two conditions are equivalent:

- (i) H is the algebraic direct sum of M_1 and M_2 .
- (*ii*) $H = M_1 + M_2$ and $M_1 \cap M_2 = \{0\}$.

Proof. Assume (*ii*) holds and let $x \in H$. Then we have $x = m_1 + m_2$ for some $m_1 \in M_1$ and $m_2 \in M_2$. Assume we also have $x = m'_1 + m'_2$ for some $m'_1 \in M_1$ and $m'_2 \in M_2$. Then we get

$$m_1 - m'_1 = m'_2 - m_2 \in M_1 \cap M_2$$
.

Since $M_1 \cap M_2 = \{0\}$, this implies that $m'_1 = m_1$ and $m'_2 = m_2$. Thus (i) holds. Conversely, assume (i) holds. It then obvious that $H = M_1 + M_2$. Consider $y \in M_1 \cap M_2$. Then we have y = y + 0 with $y \in M_1$, $0 \in M_2$, and y = 0 + y with $0 \in M_1$, $y \in M_2$. As y has a unique decomposition as such a sum of vectors, we get that y = 0. This shows that $M_1 \cap M_2 = \{0\}$, hence that (i) holds.

Corollary 3.2.2 is the key to the following fundamental result.

Theorem 3.2.5. Let M be a closed subspace of a Hilbert space H and set $M^{\perp} = \{z \in H : \langle z, y \rangle = 0 \text{ for all } y \in M\}.$

(i) If $x \in H$ and $x_M \in M$ is the best approximation to x in M, then we have $x - x_M \in M^{\perp}$ and

$$x = x_M + (x - x_M).$$

Moreover, M^{\perp} is a closed subspace of H and we have

$$H = M \oplus M^{\perp}.$$

(ii) The map $P_M : H \to H$ defined by $P_M(x) = x_M$ for $x \in H$ is called the orthogonal projection of H on M. We sometimes write Proj_M instead of P_M . It satisfies the following properties:

• P_M is linear and bounded, with $||P_M|| = 1$ if $M \neq \{0\}$.

•
$$(P_M)^2 = P_M$$

• $P_M(H) = M$ and $\ker(P_M) = M^{\perp}$.

(iii) We also have

$$(M^{\perp})^{\perp} = M$$
 and $P_{M^{\perp}} = I_H - P_M$.

Proof. (i) Let $x \in H$ and set $x^{\perp} := x - x_M$. We claim that x^{\perp} lies in M^{\perp} .

To show this, let $y \in M$ and $\varepsilon > 0$. Since $(x_M + \varepsilon y) \in M$, we get from Theorem 3.2.1 that

$$\|x^{\perp}\|^{2} = \|x - x_{M}\|^{2} \leq \|x - (x_{M} + \varepsilon y)\|^{2} = \|x^{\perp} - \varepsilon y\|^{2}$$
$$= \|x^{\perp}\|^{2} - 2\varepsilon \operatorname{Re}(\langle x^{\perp}, y \rangle) + \varepsilon^{2} \|y\|^{2},$$

which gives that

$$2 \operatorname{Re}(\langle x^{\perp}, y \rangle) \leq \varepsilon \, \|y\|^2 \, .$$

As this holds for every $\varepsilon > 0$, we obtain that $\operatorname{Re}(\langle x^{\perp}, y \rangle) \leq 0$. Applying this to $-y \in M$, we also get that $-\operatorname{Re}(\langle x^{\perp}, y \rangle) \leq 0$, i.e., $\operatorname{Re}(\langle x^{\perp}, y \rangle) \geq 0$. Thus, it follows that $\operatorname{Re}(\langle x^{\perp}, y \rangle) = 0$. If $\mathbb{F} = \mathbb{R}$, this means that $\langle x^{\perp}, y \rangle = 0$. If $\mathbb{F} = \mathbb{C}$, we also have that $iy \in M$, and this gives that

$$\operatorname{Im}(\langle x^{\perp}, y \rangle) = \operatorname{Re}(-i \langle x^{\perp}, y \rangle) = \operatorname{Re}(\langle x^{\perp}, i y \rangle) = 0.$$

Thus, $\langle x^{\perp}, y \rangle = 0$ in this case too. As this holds for every $y \in M$, the claim is proven.

Now, $x = x_M + (x - x_M)$, and $x - x_M = x^{\perp} \in M^{\perp}$. This shows the first part of (i) and that

$$H = M + M^{\perp}$$

We also have that $M \cap M^{\perp} = \{0\}$: indeed, if $y \in M \cap M^{\perp}$, then $\langle y, y \rangle = 0$, so y = 0.

Finally, it is an easy exercise to check that M^{\perp} is a closed subspace of H. Altogether, we get that $H = M \oplus M^{\perp}$, as desired.

(*ii*) We also leave it as an exercise to verify that the map P_M is linear, and prove its other properties. Using Pythagoras' identity, we get that

$$||P_M(x)||^2 = ||x_M||^2 \le ||x_M||^2 + ||x - x_M||^2 = ||x_M + (x - x_M)||^2 = ||x||^2$$

for all $x \in H$. This shows that P_M is bounded with $||P_M|| \leq 1$.

If $y \in M$, then y = y + 0 and $0 \in M^{\perp}$; this implies that $P_M(y) = y$. Hence, $||P_M(y)|| = 1$ if $y \in M$ and ||y|| = 1. This implies that $||P_M|| \ge 1$, hence that $||P_M|| = 1$, if $M \neq \{0\}$.

Moreover, let $x \in H$. Since $P_M(x) = x_M \in M$ we get that

$$(P_M)^2(x) = P_M(x_M) = x_M = P_M(x).$$

Thus $(P_M)^2 = P_M$. We also see that $P_M(H) \subseteq M$. On the other hand, if $y \in M$, then $y = P_M(y) \in P_M(H)$. Thus $M \subseteq P_M(H)$, and it follows that $P_M(H) = M$.

Finally, assume that $x \in \ker(P_M)$, i.e., $x_M = 0$. Then

$$x = x - 0 = x - x_M \in M^\perp.$$

Conversely, assume that $x \in M^{\perp}$. Since M^{\perp} is a subspace, $x_M - x = -(x - x_M) \in M^{\perp}$, hence $x_M = x + (x_M - x) \in M^{\perp}$. Altogether, this shows that $\ker(P_M) = M^{\perp}$.

(*iii*) Let $y \in M$. Then for all $z \in M^{\perp}$, we have $\langle y, z \rangle = \overline{\langle z, y \rangle} = 0$. This implies that $y \in (M^{\perp})^{\perp}$. Hence we have $M \subseteq (M^{\perp})^{\perp}$.

To show that $(M^{\perp})^{\perp} \subseteq M$, we first observe that by applying the first part of the theorem to M^{\perp} , we get that

$$H = M^{\perp} \oplus (M^{\perp})^{\perp}$$

Now let $x \in (M^{\perp})^{\perp}$, and set $x^{\perp} := x - x_M \in M^{\perp}$. Since $x_M \in M \subseteq (M^{\perp})^{\perp}$, we can write

$$x = x^{\perp} + x_M$$
, where $x^{\perp} \in M^{\perp}$ and $x_M \in (M^{\perp})^{\perp}$, and $x = 0 + x$, where $0 \in M^{\perp}$ and $x \in (M^{\perp})^{\perp}$.

By the uniqueness of decomposition in a direct sum, we get that $x = x_M$, so $x \in M$. Thus, we have shown that $(M^{\perp})^{\perp} \subseteq M$, and we can conclude that $(M^{\perp})^{\perp} = M$.

Finally, for $x \in H$, we have

$$x = (x - x_M) + x_M \,,$$

where $(x - x_M) \in M^{\perp}$ and $x_M \in M = (M^{\perp})^{\perp}$. This gives that

$$P_{M^{\perp}}(x) = x - x_M = (I_H - P_M)(x).$$

Hence, $P_{M^{\perp}} = I_H - P_M$.

Remark 3.2.6. Assume that M is finite-dimensional subspace of a Hilbert space H and that $\mathcal{B} = \{u_1, \ldots, u_n\}$ is an orthonormal basis for M. Then we know that the orthogonal projection P_M of H on M is given by

$$P_M(x) = \sum_{j=1}^n \langle x, u_j \rangle \, u_j \quad \text{for all } x \in H$$

A similar formula holds when M is only assumed to be a closed subspace of H, as we will see in the next section after having discussed orthonormal bases in Hilbert spaces.

An immediate, but noteworthy, consequence of Theorem 3.2.5 is:

Corollary 3.2.7. Let M be closed subspace of a Hilbert space H. Then M = H if and only if $M^{\perp} = \{0\}$.

In connection with the next corollary, we recall that if S is a nonempty subset of a vector space V, then Span (S) denote the subspace of V consisting of all possible finite linear combinations of vectors in S.

Corollary 3.2.8. Let S denote a nonempty subset of a Hilbert space H. Then Span(S) is dense in H if and only if $S^{\perp} = \{0\}$.

Proof. Set $M := \overline{\text{Span}(S)}$, which is a closed subspace of H. Then Span(S) is dense in H if and only if M = H. As one easily verifies that $S^{\perp} = M^{\perp}$ (cf. Exercise 3.5), the result follows from Corollary 3.2.7.

A nonempty subset S of a normed space X is sometimes called *total in* X when Span (S) is dense in X. So the corollary above says that S is total in H if and only if $S^{\perp} = \{0\}$.

Example 3.2.9. Let (X, \mathcal{A}, μ) be a measure space and $L^2 := L^2(X, \mathcal{A}, \mu)$.

Let $E \in \mathcal{A}$ and set $F := E^c \in \mathcal{A}$. If $g : X \to \mathbb{C}$ is measurable, we will say that g lives essentially on E when $\mu(\{x \in F : g(x) \neq 0\}) = 0$. Then we will let M_E denote the subset of L^2 given by

$$M_E := \left\{ [g] : g \in \mathcal{L}^2 \text{ and } g \text{ lives essentially on } E \right\}.$$

Similarly, we can define $M_F \subseteq L^2$. We claim that

$$M_F = (M_E)^{\perp}$$
 and $M_E = (M_F)^{\perp}$. (3.2.2)

To prove this, assume first that $[g] \in M_E$ and $[h] \in M_F$. Then one easily sees that $g = g \mathbf{1}_E \mu$ -a.e. and $h = h \mathbf{1}_F \mu$ -a.e. As $E \cap F = \emptyset$, we get

$$\langle [g], [h] \rangle = \int_X g \, \mathbf{1}_E \, \overline{h} \, \mathbf{1}_F \, \mathrm{d}\mu = \int_X g \, \overline{h} \, \mathbf{1}_{E \cap F} \, \mathrm{d}\mu = 0$$

Since this is true for all $[g] \in M_E$, this implies that $[h] \in (M_E)^{\perp}$. As this holds for all $[h] \in M_F$, we get that $M_F \subseteq (M_E)^{\perp}$.

To show the reverse inclusion, let $[h] \in (M_E)^{\perp}$. Then we have

$$\int_X g \,\overline{h} \, \mathrm{d}\mu = 0 \quad \text{whenever } [g] \in M_E \,.$$

In particular, since $[h \mathbf{1}_E] \in M_E$, we get

$$\int_X |h|^2 \mathbf{1}_E \, \mathrm{d}\mu = \int_X (h \, \mathbf{1}_E) \,\overline{h} \, \mathrm{d}\mu = 0$$

Since $|h|^2 \mathbf{1}_E$ is nonnegative on X, this implies that

$$\mu(\{x \in X : |h(x)|^2 \mathbf{1}_E(x) \neq 0\}) = 0.$$

As $\{x \in E : h(x) \neq 0\} = \{x \in X : |h(x)|^2 \mathbf{1}_E(x) \neq 0\}$, we get that $\mu(\{x \in E : h(x) \neq 0\}) = 0$, hence that h lives essentially on F. Thus, $[h] \in M_F$. This shows that $(M_E)^{\perp} \subseteq M_F$.

Altogether, we have proved that $M_F = (M_E)^{\perp}$. Interchanging E and F, we get that $M_E = (M_F)^{\perp}$, and the proof of (3.2.2) is finished.

Since the orthogonal complement of any subset is a closed subspace, we can conclude that M_E and M_F are closed subspaces of L^2 . Theorem 3.2.5 now gives that

$$L^2 = M_E \oplus (M_E)^{\perp} = M_E \oplus M_F$$

We note that the fact that $L^2 = M_E + M_F$ is a simple consequence of the equation

$$[f] = [f \mathbf{1}_E] + [f \mathbf{1}_F]$$
, where $[f \mathbf{1}_E] \in M_E$, $[f \mathbf{1}_F] \in M_F$,

which holds for all $[f] \in L^2$. From this equation, we now see that the orthogonal projection of L^2 on M_E (resp. M_F) is given by

$$P_{M_E}([f]) = [f \mathbf{1}_E]$$
 (resp. $P_{M_F}([f]) = [f \mathbf{1}_F]$).

3.3 Orthonormal bases in Hilbert spaces

The notion of an orthonormal basis for a finite-dimensional inner product space, which is well-known from elementary linear algebra, have a natural generalization to Hilbert spaces.

Definition 3.3.1. A nonempty subset \mathcal{B} of a Hilbert space H is called an *orthonormal basis* for H when \mathcal{B} is orthonormal and Span (\mathcal{B}) is dense in H.

Suppose H is a (nonzero) finite-dimensional Hilbert space. Then an orthonormal set \mathcal{B} in H has to be finite, so $\text{Span}(\mathcal{B})$, being finite-dimensional, is closed in H; hence, $\text{Span}(\mathcal{B})$ is dense in H if and only if $\text{Span}(\mathcal{B}) = H$. Thus we see that Definition 3.3.1 agrees with the usual one when H is

finite-dimensional. As a curiosity, we also mention that some authors like to define the empty set \emptyset to be an orthonormal basis for the trivial Hilbert space $H = \{0\}$.

Our first example is of great importance in Fourier analysis.

Example 3.3.2. Let $H = L^2([-\pi, \pi], \mathcal{A}, \mu)$, where \mathcal{A} denotes the σ -algebra of all Lebesgue measurable subsets of $[-\pi, \pi]$, and μ is the *normalized* Lebesgue measure on \mathcal{A} , that is,

$$\mu(A) := \frac{1}{2\pi} \lambda(A) \text{ for all } A \in \mathcal{A},$$

where λ denotes the Lebesgue measure on \mathbb{R} . In particular, we have $\mu([-\pi,\pi]) = 1$. For each $n \in \mathbb{Z}$, let $e_n : [-\pi,\pi] \to \mathbb{C}$ denote the continuous function given by

$$e_n(t) := e^{int}$$
 for all $t \in [-\pi, \pi]$.

As should be well-known (and is easy to check), the set

$$\mathcal{B} := \{ [e_n] : n \in \mathbb{Z} \}$$

is an orthonormal subset of H. We claim that $\text{Span}(\mathcal{B})$ is dense in H.

To show this claim, let \mathcal{T} denote the space of all (complex) trigonometrical polynomials, i.e., $\mathcal{T} := \text{Span}(\{e_n : n \in \mathbb{Z}\})$. Clearly, we have

$$\operatorname{Span}\left(\mathcal{B}\right) = \left\{\left[h\right] : h \in \mathcal{T}\right\}$$

Further, let $C_{per}([-\pi,\pi]) = \{k \in C([-\pi,\pi]) : k(-\pi) = k(\pi)\}$. We will use the fact (shown for example in Lindstrøm's book) that \mathcal{T} is dense in $C_{per}([-\pi,\pi])$ w.r.t. the uniform norm $\|\cdot\|_u$.

Let $[f] \in H$ and $\varepsilon > 0$. Using Proposition 2.1.7 we can find some $g \in C([-\pi, \pi])$ such that

$$\|[f] - [g]\|_2 < \varepsilon/3. \tag{3.3.1}$$

Further, it is not difficult to see that we can pick $k \in C_{per}([-\pi,\pi])$ such that

$$||[g] - [k]||_2 = ||g - k||_2 < \varepsilon/3.$$
(3.3.2)

Now, as mentioned above, we can find $h \in \mathcal{T}$ such that $||k - h||_u < \varepsilon/3$. Since

$$\| [k] - [h] \|_{2}^{2} = \int_{[-\pi,\pi]} |k - h|^{2} d\mu$$

$$\leq \| k - h \|_{u}^{2} \int_{[-\pi,\pi]} d\mu$$

$$= \| k - h \|_{u}^{2} \mu([-\pi,\pi])$$

$$= \| k - h \|_{u}^{2},$$

we get that

$$\|[k] - [h]\|_{2} \le \|k - h\|_{u} < \varepsilon/3.$$
(3.3.3)

Using the triangle inequality, (3.3.1), (3.3.2) and (3.3.3), we obtain that

$$\begin{split} \| [f] - [h] \|_2 &= \| [f] - [g] + [g] - [k] + [k] - [h] \|_2 \\ &\leq \| [f] - [g] \|_2 + \| [g] - [k] \|_2 + \| [k] - [h] \|_2 \\ &< \varepsilon/3 + \varepsilon/3 + \varepsilon/3 = \varepsilon \,. \end{split}$$

This shows that $[f] \in \overline{\text{Span}(\mathcal{B})}$. Hence, $\overline{\text{Span}(\mathcal{B})} = H$, as claimed above.

We can therefore conclude that $\mathcal{B} = \{[e_n] : n \in \mathbb{Z}\}$ is an orthonormal basis for H.

More generally, we may consider the L^2 -space associated to an interval [a, b] and the normalized Lebesgue measure $\mu := \frac{1}{b-a} \lambda$. Then, letting e'_n be defined for each $n \in \mathbb{Z}$ by

$$e'_n(t) = e^{int \, 2\pi/(b-a)} \quad \text{for all } t \in [a, b],$$

one may argue in a similar way as above, and conclude that $\mathcal{B}' = \{e'_n : n \in \mathbb{Z}\}$ is an orthonormal basis for this L^2 -space.

An immediate consequence of Corollary 3.2.8 is the following useful characterization of orthonormal bases:

Proposition 3.3.3. Assume that \mathcal{B} is an orthonormal subset of a Hilbert space H. Then \mathcal{B} is an orthonormal basis for H if and only if $\mathcal{B}^{\perp} = \{0\}$.

Example 3.3.4. Let X be a nonempty set. Then $\ell^2(X)$ has a natural orthonormal basis \mathcal{B} which is the analogue of the standard basis $\{e_1, \ldots, e_n\}$ for \mathbb{F}^n (which may be identified with $\ell^2(\{1, \ldots, n\})$).

Indeed, for each $x \in X$, let $e_x \in \ell^2(X)$ be defined by $e_x = \mathbf{1}_{\{x\}}$, and set

$$\mathcal{B} := \{e_x : x \in X\}.$$

Then \mathcal{B} is clearly orthonormal. Moreover, let $f \in \ell^2(X)$, $f \in \mathcal{B}^{\perp}$. For each $x \in X$, we have $\langle f, e_x \rangle = 0$. As

$$\langle f, e_x \rangle = \sum_{y \in X} f(y) \overline{e_x(y)} = \sum_{y \in \{x\}} f(y) = f(x),$$

we get that f(x) = 0 for all $x \in X$, i.e., f = 0. This shows that $\mathcal{B}^{\perp} = \{0\}$, and Proposition 3.3.3 gives that \mathcal{B} is an orthonormal basis for $\ell^2(X)$. Note that \mathcal{B} is uncountable if X is uncountable, e.g. if $X = \mathbb{R}$.

It can be shown that every (non-zero) Hilbert space has an orthonormal basis. The proof is nonconstructive as it relies on Zorn's lemma, i.e., on the axiom of choice. We will take this fact as granted.

Example 3.3.5. The *Gram-Schmidt orthonormalization prosess*, of great usefulness in the finite-dimensional case, can be generalized to cover the following situation:

Let *H* be a Hilbert space, $H \neq \{0\}$. Let $\{x_j\}_{j \in \mathbb{N}}$ be a sequence of vectors in $H \setminus \{0\}$ and set $S := \{x_j : j \in \mathbb{N}\}$. Assume that Span (S) is dense in *H*.

We remark that such a sequence exists whenever H is finite-dimensional (since repetitions are allowed in a sequence). More generally, it exists whenever H is *separable*, i.e., whenever H contains a countable dense subset, cf. Exercise 3.10.

For each $n \in \mathbb{N}$, set $M_n := \text{Span}(\{x_1, \ldots, x_n\})$. We note that for each n we have $M_n \subseteq M_{n+1}$. Moreover, $\text{Span}(S) = \bigcup_{n \in \mathbb{N}} M_n$.

Proceeding inductively, we can construct an orthonormal basis \mathcal{B}_n for each M_n as follows:

- i) We set $\mathcal{B}_1 := \left\{ \frac{1}{\|x_1\|} x_1 \right\}$. Clearly, \mathcal{B}_1 is an orthonormal basis for M_1 .
- ii) Let $n \in \mathbb{N}$ and assume that we have constructed an orthonormal basis \mathcal{B}_n for M_n .

If $x_{n+1} \in M_n$, then set $\mathcal{B}_{n+1} := \mathcal{B}_n$. Otherwise, set

$$y_{n+1} := x_{n+1} - \operatorname{Proj}_{M_n}(x_{n+1}) \text{ and } \mathcal{B}_{n+1} := \mathcal{B}_n \cup \left\{ \frac{1}{\|y_{n+1}\|} \, y_{n+1} \right\}.$$

It follows readily that \mathcal{B}_{n+1} is an orthonormal basis for M_{n+1} .

Set now $\mathcal{B} := \bigcup_{n \in \mathbb{N}} \mathcal{B}_n$. Then \mathcal{B} is orthonormal, and $\text{Span}(\mathcal{B}) = \text{Span}(S)$, so

$$\overline{\operatorname{Span}(\mathcal{B})} = \overline{\operatorname{Span}(S)} = H$$
.

Hence, \mathcal{B} is an orthonormal basis for H.

We observe that since each \mathcal{B}_n is finite, \mathcal{B} is countable. Conversely, if H has a countable orthonormal basis, then it can be shown that H is separable (cf. Exercise 3.10).

When H is a nontrivial finite-dimensional inner product space, and $\mathcal{B} = \{u_1, \ldots, u_n\}$ is an orthonormal basis for H, we know that every $x \in H$ has a Fourier expansion w.r.t. \mathcal{B} , i.e., we have

$$x = \sum_{j=1}^{n} \langle x, u_j \rangle \, u_j \, .$$

As we will soon see, a similar expansion also holds in any infinite dimensional Hilbert space.

We will use the following notation. If J is a nonempty set (possibly uncountable), and $j \mapsto t_j$ is a function from J into $[0, \infty)$, we set

$$\sum_{j \in J} t_j := \sup \left\{ \sum_{j \in F} t_j : F \subseteq J, F \text{ is finite and nonempty} \right\} \in [0, \infty].$$

Equivalently, $\sum_{j \in J} t_j$ is the integral of the nonnegative function $j \mapsto t_j$ w.r.t. the counting measure on $\mathcal{P}(J)$ (= the σ -algebra of all subsets of J).

We first note that Bessel's inequality holds for any orthonormal set:

Lemma 3.3.6. Assume that \mathcal{B} is an orthonormal set in an inner product space H, and let $x \in H$. Then

$$\sum_{u \in \mathcal{B}} \left| \langle x, u \rangle \right|^2 \le \|x\|^2.$$

and the set $\mathcal{B}_x := \{ u \in \mathcal{B} : \langle x, u \rangle \neq 0 \}$ is countable.

Proof. Let F be a nonempty finite subset of \mathcal{B} . As F is orthonormal, Bessel's inequality for F gives that

$$\sum_{u \in F} \left| \langle x, u \rangle \right|^2 \le \|x\|^2.$$

Thus we get that

$$\sup\left\{\sum_{u\in F} \left|\langle x,u\rangle\right|^2: F\subseteq \mathcal{B}, \ F \text{ is finite and nonempty}\right\} \le \|x\|^2,$$

which proves the first assertion.

Further, this implies that the set $\mathcal{B}_{x,n} := \{u \in \mathcal{B} : |\langle x, u \rangle|^2 \ge 1/n\}$ is finite for every $n \in \mathbb{N}$. Hence, $\mathcal{B}_x = \bigcup_{n \in \mathbb{N}} \mathcal{B}_{x,n}$ is countable.

On Hilbert spaces and bounded linear operators

The next lemma will be useful at several occasions.

Lemma 3.3.7. Assume $\{u_j : j \in \mathbb{N}\}$ is a countably infinite orthonormal subset of a Hilbert space H, and let $\{c_j\}_{j\in\mathbb{N}}$ be a sequence in \mathbb{F} satisfying that

$$\sum_{j=1}^{\infty} |c_j|^2 < \infty$$
 .

Then the series $\sum_{j=1}^{\infty} c_j u_j$ converges to some $y \in H$, and we have that

$$\langle y, u_k \rangle = c_k$$
 for every $k \in \mathbb{N}$.

Proof. A similar result is shown in Lindstrøm's book, but we sketch the argument for the ease of the reader. For each $n \in \mathbb{N}$, set $y_n = \sum_{j=1}^n c_j u_j$. Then, for any m > n, Pythagoras' identity gives that

$$||y_m - y_n||^2 = \sum_{j=n+1}^m ||c_j u_j||^2 = \sum_{j=n+1}^m |c_j|^2.$$

Using the assumption, the sum above can be made as small as we want by choosing m and n large enough. Thus the sequence $\{y_n\}_{n\in\mathbb{N}}$ is Cauchy in H, so it converges to some $y \in H$, i.e., we have $y = \sum_{j=1}^{\infty} c_j u_j$. For each $k \in \mathbb{N}$, continuity and linearity of the inner product in the first variable gives then that $\langle y, u_k \rangle = \sum_{j=1}^{\infty} c_j \langle u_j, u_k \rangle = c_k$.

Theorem 3.3.8. Let H be a Hilbert space, $H \neq \{0\}$, and let \mathcal{B} be an orthonormal subset of H. Then the following conditions are equivalent:

- (a) \mathcal{B} is an orthonormal basis for H.
- (b) Every $x \in H \setminus \{0\}$ has a Fourier expansion

$$x = \sum_{u \in \mathcal{B}_x} \langle x, u \rangle \, u \tag{3.3.4}$$

where $\mathcal{B}_x = \{u \in \mathcal{B} : \langle x, u \rangle \neq 0\}$ is countable (cf. Lemma 3.3.6) and nonempty.

In the case where \mathcal{B}_x is countably infinite, we mean by (3.3.4) that the following holds: if $\mathcal{B}_x = \{u_j : j \in \mathbb{N}\}$ is any enumeration of the distinct elements of \mathcal{B}_x , then we have

$$\lim_{n \to \infty} \left\| x - \sum_{j=1}^{n} \langle x, u_j \rangle \, u_j \right\| = 0 \,, \, i.e., \, x = \sum_{j=1}^{\infty} \langle x, u_j \rangle \, u_j \,.$$
(c) For every $x \in H$ we have $\|x\|^2 = \sum_{u \in \mathcal{B}} \left| \langle x, u \rangle \right|^2$.

The formula in (c) is called Parseval's identity.

Proof. (a) \Rightarrow (b): Assume that \mathcal{B} is an orthonormal basis for H and let $x \in H \setminus \{0\}$.

We first observe that $\mathcal{B}_x \neq \emptyset$. Indeed, suppose that $\mathcal{B}_x = \emptyset$. This means that $x \in \mathcal{B}^{\perp}$. But $\mathcal{B}^{\perp} = \{0\}$ by Proposition 3.3.3, so x = 0, a contradiction.

We only consider the case where \mathcal{B}_x is countably infinite. (The case where \mathcal{B}_x is finite is much easier and left to the reader). Let $\{u_j : j \in \mathbb{N}\}$ be an enumeration of the distinct elements of \mathcal{B}_x . Since \mathcal{B}_x is orthonormal, Bessel's inequality gives that

$$\sum_{j=1}^{\infty} \left| \langle x, u_j \rangle \right|^2 \le ||x||^2.$$

Applying Lemma 3.3.7 with $c_j = \langle x, u_j \rangle$ for every $j \in \mathbb{N}$, we get that the series $\sum_{j=1}^{\infty} \langle x, u_j \rangle u_j$ converges to some $y \in H$, which satisfies that

$$\langle y, u_k \rangle = c_k = \langle x, u_k \rangle$$
 for every $k \in \mathbb{N}$.

Moreover, if $u \in \mathcal{B} \setminus \mathcal{B}_x$, we get that

$$\langle y, u \rangle = \sum_{j=1}^{\infty} \langle x, u_j \rangle \langle u_j, u \rangle = 0 = \langle x, u \rangle.$$

It follows that $x - y \in \mathcal{B}^{\perp} = \{0\}$, hence that x = y. This shows that the assertion in (b) holds in this case.

 $(b) \Rightarrow (c)$: Assume (b) holds, and let $x \in H \setminus \{0\}$. Again we consider the more difficult case where \mathcal{B}_x is countably infinite, so we may write $\mathcal{B}_x = \{u_j : j \in \mathbb{N}\}$ as above. By continuity of the norm and Pythagoras' identity, we get

$$||x||^2 = \sum_{j=1}^{\infty} \left| \langle x, u_j \rangle \right|^2.$$

Hence, given $\varepsilon > 0$, we can find $n \in \mathbb{N}$ such that $||x||^2 - \sum_{j=1}^n |\langle x, u_j \rangle|^2 < \varepsilon$, giving

$$||x||^2 - \varepsilon < \sum_{j=1}^n |\langle x, u_j \rangle|^2 \le \sum_{u \in \mathcal{B}} |\langle x, u \rangle|^2.$$

Since this holds for every $\varepsilon > 0$, we get that $||x||^2 \leq \sum_{u \in \mathcal{B}} |\langle x, u \rangle|^2$. Combining this inequality with Lemma 3.3.6, we see that (c) holds.

(c) \Rightarrow (a): Assume $||x||^2 = \sum_{u \in \mathcal{B}} |\langle x, u \rangle|^2$ for every $x \in H$. If $x \in \mathcal{B}^{\perp}$, i.e., $\langle x, u \rangle = 0$ for every $u \in \mathcal{B}$, then we get $||x||^2 = 0$, so x = 0. Hence, $\mathcal{B}^{\perp} = \{0\}$, and Proposition 3.3.3 gives that (a) holds.

Remark 3.3.9. The Fourier expansion of x in Theorem 3.3.8 (b) can be written in the form

$$x = \sum_{u \in \mathcal{B}} \langle x, u \rangle \, u \tag{3.3.5}$$

if one takes care of giving a meaning to convergence of generalized sums in normed spaces. We discuss this in Exercise 3.16. In these notes, we will sometimes use (3.3.5) as a short form of the Fourier expansion of x given by (3.3.4).

Example 3.3.10. Let M be a closed subspace of a Hilbert space H and assume that we have found an orthonormal basis C for M. Then we can find a formula for the orthogonal projection P_M of H on M:

Let $x \in H$. If $x \in M^{\perp}$, then $P_M(x) = 0$, so we can assume $x \in H \setminus M^{\perp}$. Since \mathcal{C} is orthonormal in H, we know that $\mathcal{C}_x := \{v \in \mathcal{C} : \langle x, v \rangle \neq 0\}$ is countable. Set $x_M := P_M(x) \in M$ and $x^{\perp} := x - x_M \in M^{\perp}$, and note that $x_M \neq 0$. Now, for each $v \in \mathcal{C}$, we have

$$\langle x, v \rangle = \langle x_M, v \rangle + \langle x^{\perp}, v \rangle = \langle x_M, v \rangle.$$

Hence, $C_x = C_{x_M}$. Moreover, applying Theorem 3.3.8 to $M, x_M \in M \setminus \{0\}$ and C, we get that

$$x_M = \sum_{v \in \mathcal{C}_{x_M}} \langle x_M, v \rangle v$$

Using our previous observations, this formula can be rewritten as

$$P_M(x) = \sum_{v \in \mathcal{C}_x} \langle x, v \rangle v$$

It generalizes the usual formula for $P_M(x)$ when M is finite-dimensional.

An immediate consequence of Theorem 3.3.8 is the following:

Corollary 3.3.11. Assume a Hilbert space H contains a countably infinite orthonormal subset \mathcal{B} , enumerated as $\mathcal{B} = \{v_k : k \in \mathbb{N}\}$. Then \mathcal{B} is an orthonormal basis for H if and only if

$$x = \sum_{k=1}^{\infty} \langle x, v_k \rangle \, v_k$$

for all $x \in H$, if and only if

$$||x||^{2} = \sum_{k=1}^{\infty} \left| \langle x, v_{k} \rangle \right|^{2}$$

for all $x \in H$.

Example 3.3.12. Let $\mathcal{B} = \{[e_n] : n \in \mathbb{Z}\}$ denote the standard orthonormal basis for $H = L^2([-\pi, \pi])$ described in Example 3.3.2. For $[f] \in H$ and $n \in \mathbb{Z}$ it is common to set

$$\widehat{[f]}(n) := \langle [f], [e_n] \rangle = \frac{1}{2\pi} \int_{[-\pi,\pi]} f(t) e^{-int} d\lambda(t) \,,$$

which is called the *Fourier coefficient* of [f] at n.

In fact, it is usual to write f instead of [f], having in mind that one then identifies functions which agree μ -a.e. Hence, the Fourier coefficient of f at n is denoted by $\hat{f}(n)$, and the Fourier expansion of f w.r.t. \mathcal{B} is then written as

$$f = \sum_{n \in \mathbb{Z}} \widehat{f}(n) e_n \,,$$

meaning that

$$f = \lim_{m \to \infty} \sum_{n = -m}^{m} \widehat{f}(n) e_n \quad (\text{w.r.t. } \| \cdot \|_2).$$

This follows from Corollary 3.3.11 by enumerating \mathcal{B} as e_0 , e_{-1} , e_1 , e_{-2} , e_2 , etc. Similarly, we have

$$||f||_{2}^{2} = \sum_{n \in \mathbb{Z}} \left| \widehat{f}(n) \right|^{2}.$$

3.4 Adjoint operators

Let X be a normed space (over \mathbb{F}). We recall that the dual space X^* consists of the bounded linear functionals on X (with values in \mathbb{F}), and that X^* is a Banach space w.r.t. the norm $\|\varphi\| := \sup\{|\varphi(x)| : x \in X_1\}, \varphi \in X^*$. An important goal of functional analysis is to gain new insight by exploiting the interplay between a space and its dual. This is particularly successful when X is a Hilbert space because the dual space may then be identified in a natural way with the space itself.

Theorem 3.4.1. Let H be a Hilbert space (over \mathbb{F}). For each $y \in H$, define $\varphi_y : H \to \mathbb{F}$ by

$$\varphi_y(x) := \langle x, y \rangle \quad \text{for all } x \in H.$$

Then $\varphi_y \in H^*$ for all $y \in H$.

Moreover, the map $y \mapsto \varphi_y$ is a bijection from H onto H^* , which is isometric, and conjugate-linear in the sense that

$$\varphi_{\lambda_1 y_1 + \lambda_2 y_2} = \lambda_1 \, \varphi_{y_1} + \lambda_2 \, \varphi_{y_2}$$

for all $\lambda_1, \lambda_2 \in \mathbb{F}$ and all $y_1, y_2 \in H$.

Proof. Let $y \in H$. Then the map φ_y is clearly linear. Moreover, for all $x \in H$, the Cauchy-Schwarz inequality gives that

$$|\varphi_y(x)| = |\langle x, y \rangle| \le ||x|| ||y||.$$

Hence, φ_y is bounded, with $\|\varphi_y\| \leq \|y\|$. If $y \neq 0$, then

$$\left|\varphi_{y}\left(\frac{1}{\|y\|}y\right)\right| = \frac{1}{\|y\|} \langle y, y \rangle = \|y\|,$$

so $\|\varphi_y\| \ge \|y\|$. Thus, $\|\varphi_y\| = \|y\|$.

This shows that the map $y \mapsto \varphi_y$ is an isometry from H into H^* . In particular, it is injective. To show that it is surjective, let $\varphi \in H^*$. If $\varphi = 0$, then we have $\varphi = \varphi_0$. So assume $\varphi \neq 0$ and set $M := \ker \varphi$. Then M is a closed subspace of H such that $M \neq H$. By Corollary 3.2.7, $M^{\perp} \neq \{0\}$, so we can pick $z \in M^{\perp}$ such that ||z|| = 1, and set

$$y := \overline{\varphi(z)} \, z \in H \, .$$

We claim that $\varphi = \varphi_y$. Indeed, let $x \in H$ and set $m := \varphi(x) z - \varphi(z) x \in H$. Then we have

$$\varphi(m) = \varphi(x)\,\varphi(z) - \varphi(z)\,\varphi(x) = 0\,,$$

so $m \in M$. As $z \in M^{\perp}$, we get $\langle m, z \rangle = 0$, i.e.,

$$\langle \varphi(x) \, z, z \rangle = \langle \varphi(z) \, x, z \rangle.$$

Hence

$$\varphi(x) = \varphi(x) ||z||^2 = \langle \varphi(x) z, z \rangle = \langle \varphi(z) x, z \rangle = \varphi(z) \langle x, z \rangle$$
$$= \langle x, \overline{\varphi(z)} z \rangle = \langle x, y \rangle = \varphi_y(x) .$$

This shows the claim that $\varphi = \varphi_y$, hence that the map $y \mapsto \varphi_y$ is surjective.

Altogether, we have shown that this map is an isometric bijection from H onto H^* , as desired.

The final assertion is an obvious consequence of the conjugate-linearity of the inner product in the second variable.

This theorem, which is one among a diversity of results being called *the Riesz representation theorem*, has several useful consequences. Our main application in these notes will be to use it to associate an adjoint operator to every bounded operator on a Hilbert space. Some people like to think of the adjoint as a kind of twin (or as a kind of shadow), which happens to coincide with the original operator in many cases of interest. **Theorem 3.4.2.** Let H be a Hilbert space (over \mathbb{F}). For each $T \in \mathcal{B}(H)$, there is a unique operator $T^* \in \mathcal{B}(H)$, called the adjoint of T, satisfying

$$\langle T(x), y \rangle = \langle x, T^*(y) \rangle \tag{3.4.1}$$

for all $x, y \in H$.

The *-operation on $\mathcal{B}(H)$, $T \mapsto T^*$, enjoys the following properties: For all $S, T \in \mathcal{B}(H)$ and all $\alpha, \beta \in \mathbb{F}$, we have

• i) $(\alpha S + \beta T)^* = \overline{\alpha} S^* + \overline{\beta} T^*;$ ii) $(ST)^* = T^*S^*;$ iii) $(T^*)^* = T;$

• *iv*)
$$||T^*|| = ||T||$$
; *v*) $||T^*T|| = ||T||^2$.

Remark 3.4.3. If H and K are Hilbert spaces (over the same \mathbb{F}), then one may associate to each $T \in \mathcal{B}(H, K)$ a unique adjoint operator $T^* \in \mathcal{B}(K, H)$ satisfying (3.4.1) for all $x \in H$ and all $y \in K$, and enjoying similar properties. We leave this as an exercise.

Proof of Theorem 3.4.2. Let $T \in \mathcal{B}(H)$ and consider $y \in H$. Using the linearity of T and the linearity of the inner product in the first variable, we get that the map $\varphi : H \to \mathbb{F}$ defined by

$$\varphi(x) := \langle T(x), y \rangle$$
 for all $x \in H$,

is a linear functional on H. Moreover, as we have

$$|\langle T(x), y \rangle| \le ||T(x)|| ||y|| \le ||T|| ||x|| ||y||$$

for all $x \in H$, φ is bounded with $\|\varphi\| \leq \|T\| \|y\|$. Hence, $\varphi \in H^*$, and Theorem 3.4.1 gives that there exists a unique vector in H, that we denote by $T^*(y)$, such that $\varphi = \varphi_{T^*(y)}$, i.e., such that

$$\langle T(x), y \rangle = \langle x, T^*(y) \rangle \tag{3.4.2}$$

for all $x \in H$. This theorem also gives that

$$||T^*(y)|| = ||\varphi_{T^*(y)}|| = ||\varphi|| \le ||T|| ||y||.$$
(3.4.3)

As what we have done above holds for every $y \in H$, we obtain a map $T^*: H \to H$ which sends each $y \in H$ to $T^*(y) \in H$. In view of (3.4.2), it is clear that (3.4.1) holds for all $x, y \in H$.

To show that T^* is linear, let $y, y' \in H$ and $\alpha \in \mathbb{F}$. Then, for all $x \in H$, we have

$$\langle x, T^*(\alpha \, y + y') \rangle = \langle T(x), \alpha \, y + y' \rangle = \overline{\alpha} \langle T(x), y \rangle + \langle T(x), y' \rangle = \overline{\alpha} \langle x, T^*(y) \rangle + \langle x, T^*(y') \rangle = \langle x, \alpha \, T^*(y) + T^*(y') \rangle .$$

This implies that $T^*(\alpha y + y') = \alpha T^*(y) + T^*(y')$, as desired.

Next, from (3.4.3), we see that T^* is bounded with $||T^*|| \leq ||T||$. To show the asserted uniqueness property of T^* , assume that $S \in \mathcal{B}(H)$ satisfies the same property as T^* , i.e.,

$$\langle T(x), y \rangle = \langle x, S(y) \rangle$$
 for all $x, y \in H$.

Let $y \in H$. Then, for all $x \in H$, we get

$$\langle x, S(y) \rangle = \langle T(x), y \rangle = \langle x, T^*(y) \rangle$$

This implies that $S(y) = T^*(y)$. Thus, $S = T^*$.

We leave the proof of properties i) and ii) as an exercise. To show the other properties, let $T \in \mathcal{B}(H)$. Then, for each $y \in H$, using equation (3.4.1) for T^* instead of T, we get that, for all $x \in H$, we have

$$\langle x, (T^*)^*(y) \rangle = \langle T^*(x), y \rangle = \overline{\langle y, T^*(x) \rangle} = \overline{\langle T(y), x \rangle} = \langle x, T(y) \rangle$$

This implies that $(T^*)^*(y) = T(y)$. Thus, $(T^*)^* = T$, i.e., *iii*) holds

Now, we have seen that $||T^*|| \leq ||T||$ holds for all $T \in \mathcal{B}(H)$. Thus we get

.

$$||T|| = ||(T^*)^*|| \le ||T^*|| \le ||T||$$

Hence $||T^*|| = ||T||$, i.e., *iv*) holds.

Further, using iv), we get $||T^*T|| \leq ||T^*|| ||T|| = ||T||^2$. On the other hand, for every $x \in H$, we have

$$\|T(x)\|^{2} = \langle T(x), T(x) \rangle = \langle x, T^{*}(T(x)) \rangle$$
$$= \left| \langle x, (T^{*}T)(x) \rangle \right| \leq \|x\| \| (T^{*}T)(x) \|$$
$$\leq \|T^{*}T\| \|x\|^{2}.$$

This implies that $||T||^2 \leq ||T^*T||$. Hence we get $||T||^2 = ||T^*T||$, i.e., v) holds.

Example 3.4.4. Consider $H = \mathbb{F}^n$ for some $n \in \mathbb{N}$ with its usual inner product, and $T \in \mathcal{B}(H)$. Let A denote the standard matrix of T. Then the standard matrix of T^* is $A^* := \overline{A}^t$.

Here, \overline{A} denotes the matrix obtained by conjugating every coefficient of A, while B^t denotes the transpose of a matrix B. Of course, if $\mathbb{F} = \mathbb{R}$, then we just get $A^* = A^t$.

Alternatively, we can formulate our assertion above by saying that if $T_A \in \mathcal{B}(H)$ denotes the operator given by multiplication with a matrix $A \in M_n(\mathbb{F})$, then we have

$$(T_A)^* = T_{A^*} \, .$$

To prove this, let $x, y \in H$. Recall that $\langle x, y \rangle = x^t \overline{y}$. So we get

$$\langle T_A(x), y \rangle = (Ax)^t \,\overline{y} = x^t A^t \,\overline{y} = x^t \,\overline{A^* y} = \langle x, T_{A^*}(y) \rangle.$$

Since this holds for all $x, y \in H$, this implies that $(T_A)^* = T_{A^*}$, as asserted.

More generally, if H is a nontrivial finite-dimensional Hilbert space, \mathcal{B} is an orthonormal basis for H, and $[T]_{\mathcal{B}}$ is the matrix of T relative to \mathcal{B} , then we have

$$[T^*]_{\mathcal{B}} = \left([T]_{\mathcal{B}} \right)^*.$$

The verification of this fact is left as an exercise. (One may argue in a similar way as in Example 3.4.6). $\hfill \square$

Example 3.4.5. Assume that $\{\lambda_j\}_{j\in\mathbb{N}}$ is bounded sequence in \mathbb{F} , and set $M := \sup_{j\in\mathbb{N}} |\lambda_j|$. If $\xi \in \ell^2$, then we have

$$\sum_{j=1}^{\infty} |\lambda_j \xi(n)|^2 \le M^2 \sum_{j=1}^{\infty} |\xi(j)|^2 = M^2 ||\xi||_2^2 < \infty$$

Thus, the map $D\xi : \mathbb{N} \to \mathbb{F}$ defined by $(D\xi)(j) = \lambda_j \xi(j)$ for every $n \in \mathbb{N}$ belongs to ℓ^2 and satisfies that

$$||D\xi||_2 \le M ||\xi||_2 \quad \text{for every } \xi \in \ell^2.$$

It follows that the map $D: \ell^2 \to \ell^2$, given by $D(\xi) = D\xi$ for each $\xi \in \ell^2$, is linear and bounded. Note that if $\{e_j\}_{j\in\mathbb{N}}$ denotes the standard orthonormal basis of ℓ^2 (cf. Exercise 3.3.4), we have that $D(e_j) = \lambda_j e_j$, hence that $\|D(e_j)\| = |\lambda_j|$ for every $j \in \mathbb{N}$. This implies that $\|D\| \ge |\lambda_j|$ for every $j \in \mathbb{N}$, hence that $\|D\| \ge M$. Altogether we get that $\|D\| = M = \sup_{j\in\mathbb{N}} |\lambda_j|$. The operator D is called the *standard diagonal operator on* ℓ^2 associated to $\{\lambda_j\}_{j\in\mathbb{N}}$.

3. On Hilbert spaces and bounded linear operators

It is easy to determine D^* : indeed, for $\xi, \eta \in \ell^2$, we have

$$\langle D(\xi),\eta\rangle = \sum_{j\in\mathbb{N}}\lambda_j\xi(j)\overline{\eta(j)} = \sum_{j\in\mathbb{N}}\xi(j)\overline{\lambda_j}\overline{\eta(j)} = \langle \xi, D^*(\eta)\rangle,$$

where $D^* \in \mathcal{B}(\ell^2)$ is the standard diagonal operator on ℓ^2 associated to the sequence $\{\overline{\lambda_j}\}_{j\in\mathbb{N}}$.

Example 3.4.6. Assume a Hilbert space H has a countably infinite orthonormal basis, enumerated as $\mathcal{B} = \{u_j\}_{j \in \mathbb{N}}$. Let $T \in \mathcal{B}(H)$. For each $(j,k) \in \mathbb{N} \times \mathbb{N}$, set

$$A(j,k) := \langle T(u_k), u_j \rangle \in \mathbb{F}.$$

We may think of the map $A : \mathbb{N} \times \mathbb{N} \to \mathbb{F}$, which sends each (j, k) to A(j, k), as the (infinite) matrix of T (w.r.t. \mathcal{B}) since, for each $k \in \mathbb{N}$, we have

$$T(u_k) = \sum_{j=1}^{\infty} \langle T(u_k), u_j \rangle \, u_j = \sum_{j=1}^{\infty} A(j,k) \, u_j \,. \tag{3.4.4}$$

Now, as every $x \in H$ has a Fourier expansion w.r.t. \mathcal{B} , it is clear that T is uniquely determined as a bounded operator on H by its values on \mathcal{B} . Thus we see from (3.4.4) that T is uniquely determined by its matrix A.

As $T^* \in \mathcal{B}(H)$ and

$$\langle T^*(u_k), u_j \rangle = \langle u_k, T(u_j) \rangle = \overline{\langle T(u_j), u_k \rangle} = \overline{A(k, j)}$$

we can conclude that the matrix of T^* w.r.t. \mathcal{B} is A^* , where

$$A^*(j,k) := \overline{A(k,j)} \,.$$

Thus we get that, for all $k \in \mathbb{N}$, we have

$$T^*(u_k) = \sum_{j=1}^{\infty} A^*(j,k) \, u_j = \sum_{j=1}^{\infty} \overline{A(k,j)} \, u_j \, .$$

From (3.4.4) and Parseval's identity, we also get that

$$\sum_{j=1}^{\infty} |A(j,k)|^2 = ||T(u_k)||^2 \le ||T||^2 < \infty$$

for each $k \in \mathbb{N}$, so the ℓ^2 -norms of the column vectors of A are uniformly bounded. However, such a condition on the column vectors of a given infinite matrix A is not sufficient in general to ensure that A is the matrix of some operator in $\mathcal{B}(H)$. There are some known conditions guaranteeing this, but we will only look below at two cases below where one can argue directly. **Example 3.4.7.** We will show how we can construct diagonal operators in certain Hilbert spaces, generalizing what we did in Example 3.4.5 in the case of ℓ^2 . Let H be a Hilbert space having a countably infinite orthonormal basis, enumerated as $\mathcal{B} = \{u_j\}_{j \in \mathbb{N}}$. Let $\{\lambda_j\}_{j \in \mathbb{N}}$ be a bounded sequence in \mathbb{F} , so that

$$M := \sup\{|\lambda_j| : j \in \mathbb{N}\} < \infty.$$

We claim that there exists an operator $D \in \mathcal{B}(H)$ satisfying that

$$D(u_k) = \lambda_k \, u_k \quad \text{for each } k \in \mathbb{N}. \tag{3.4.5}$$

Indeed, consider $x \in H$. Then Parseval's identity gives that

$$\sum_{j=1}^{\infty} \left| \lambda_j \left\langle x, u_j \right\rangle \right|^2 \le M^2 \sum_{j=1}^{\infty} \left| \left\langle x, u_j \right\rangle \right|^2 = M^2 \|x\|^2 < \infty.$$

Hence, Lemma 3.3.7 gives that the vector

$$D(x) := \sum_{j=1}^{\infty} \lambda_j \langle x, u_j \rangle \, u_j$$

satisfies that $\langle D(x), u_j \rangle = \lambda_j \langle x, u_j \rangle$ for each $j \in \mathbb{N}$. Thus, using Parseval's identity again, we get that

$$||D(x)||^2 = \sum_{j=1}^{\infty} |\lambda_j \langle x, u_j \rangle|^2 \le M^2 ||x||^2$$

It follows now readily that the map $x \mapsto D(x)$ gives an operator D in $\mathcal{B}(H)$ satisfying (3.4.5) and $||D|| \leq M$. Since $||D|| \geq ||D(u_k)|| = |\lambda_k|$ for all $k \in \mathbb{N}$, we also have that $||D|| \geq M$. Hence, ||D|| = M.

It is quite obvious that the matrix of D (w.r.t. \mathcal{B}) is the diagonal (infinite) matrix Λ defined for each $(j, k) \in \mathbb{N}$ by

$$\Lambda(j,k) = \begin{cases} \lambda_j & \text{if } j = k, \\ 0 & \text{otherwise.} \end{cases}$$

The operator D is called the diagonal operator associated to $\{\lambda_j\}_{j\in\mathbb{N}}$ (w.r.t. \mathcal{B}).

From our discussion in Example 3.4.6, we get that the matrix of D^* is Λ^* . Thus we have $D^*(u_k) = \overline{\lambda_k} u_k$ for all $k \in \mathbb{N}$, so D^* is the diagonal operator associated to $\{\overline{\lambda_j}\}_{j\in\mathbb{N}}$ (w.r.t. \mathcal{B}).

Note that if $H = \ell^2$ and we let $\mathcal{B} = \{e_j\}_{j \in \mathbb{N}}$ be its standard orthonormal basis, then the diagonal operator D associated to a bounded sequence $\{\lambda_j\}_{j \in \mathbb{N}}$ (w.r.t. \mathcal{B}) coincides with the standard diagonal operator from Example 3.4.5.

Example 3.4.8. Let again H be a Hilbert space having a countably infinite orthonormal basis, enumerated as $\mathcal{B} = \{u_j\}_{j \in \mathbb{N}}$. We will now argue that there exists an operator $S \in \mathcal{B}(H)$, called the right shift operator on H (w.r.t. \mathcal{B}), satisfying that

$$S(u_k) = u_{k+1} \quad \text{for all } k \in \mathbb{N}. \tag{3.4.6}$$

Since $\sum_{n=2}^{\infty} |\langle x, u_{n-1} \rangle|^2 = \sum_{j=1}^{\infty} |\langle x, u_j \rangle|^2 = ||x||^2 < \infty$ for all $x \in H$, we may use Lemma 3.3.7 to define a map $S: H \to H$ by

$$S(x) = \sum_{n=2}^{\infty} \langle x, u_{n-1} \rangle u_n ,$$

which is then a linear isometry satisfying that $S(u_{n-1}) = u_n$ for all $n \ge 2$, i.e., such that (3.4.6) holds. The matrix of S (w.r.t. \mathcal{B}) is the (infinite) matrix σ given by

$$\sigma(j,k) = \begin{cases} 1 & \text{if } j = k+1, \\ 0 & \text{otherwise.} \end{cases}$$

for each $(j,k) \in \mathbb{N}$. Thus, the matrix of S^* (w.r.t. \mathcal{B}) is the matrix σ^* given by

$$\sigma^*(j,k) = \overline{\sigma(k,j)} = \begin{cases} 1 & \text{if } k = j+1, \\ 0 & \text{otherwise.} \end{cases}$$

for all $j, k \in \mathbb{N}$, so we get that

$$S^*(u_k) = \sum_{j=1}^{\infty} \sigma^*(j,k) \, u_j = \begin{cases} 0 & \text{if } k = 1, \\ u_{k-1} & \text{if } k \ge 2. \end{cases}$$

The operator S^* is called the left shift operator on H (w.r.t. \mathcal{B}). We note that S^* is not isometric, in fact not even injective, because $S^*(u_1) = 0$. \Box

Example 3.4.9. Let (X, \mathcal{A}, μ) be a measure space. Set $\mathcal{L}^{\infty} := \mathcal{L}^{\infty}(X, \mathcal{A}, \mu)$ and $H := L^2(X, \mathcal{A}, \mu)$. To each $f \in \mathcal{L}^{\infty}$ we may associate a "multiplication" operator $M_f \in \mathcal{B}(H)$ given by

$$M_f([g]) = [fg] \quad \text{for all } [g] \in H.$$

Indeed, this follows readily from Proposition 2.2.4 (with q = 2). Now, for all $[g], [h] \in H$, we have

$$\langle M_f([g]), [h] \rangle = \int_X fg \,\overline{h} \, \mathrm{d}\mu = \int_X g \,\overline{\overline{f} \, h} \, \mathrm{d}\mu = \langle [g], M_{\overline{f}}([h]) \rangle \,.$$

This implies that $(M_f)^* = M_{\overline{f}}$.

Example 3.4.10. Let $K : [a, b] \times [a, b] \to \mathbb{C}$ be a continuous function and let T_K denote the associated integral operator on $H = L^2([a, b])$, which is determined on C([a, b]) by

$$[T_K(f)](s) = \int_a^b K(s,t) f(t) dt \quad \text{for all } s \in [a,b],$$

cf. Example 2.1.9 and Exercise 2.10. We leave it as Exercise 3.18 to check that $(T_K)^* = T_{K^*}$, where $K^*(s,t) := \overline{K(t,s)}$ for all $s, t \in [a,b]$.

Example 3.4.11. Let $v, w \in H$ and consider the linear operator $T_{v,w}$: $H \to H$ defined by

$$T_{v,w}(x) := \langle x, v \rangle w$$
 for all $x \in H$.

Note that $T_{v,w}$ has rank one, i.e., its range has dimension one, if $v, w \in H \setminus \{0\}$. It is an instructive exercise (cf. Exercise 3.20) to check that $T_{v,w}$ is bounded, with norm $||T_{v,w}|| = ||v|| ||w||$, and satisfies that $(T_{v,w})^* = T_{w,v}$. A bit more challenging is to use this to show that any finite-rank operator $T \in \mathcal{B}(H)$ may be written as a finite sum of rank one operators in $\mathcal{B}(H)$, and that T^* also has finite-rank. \Box

As an illustration that the adjoint operator contains valuable information about the original operator, we include a proposition showing the connection between the fundamental subspaces associated to these operators.

Proposition 3.4.12. Let H be a Hilbert space (over \mathbb{F}) and let $T \in \mathcal{B}(H)$. Then we have:

- (a) $\ker(T) = T^*(H)^{\perp}$ and $\ker(T^*) = T(H)^{\perp}$.
- (b) $\overline{T(H)} = \ker(T^*)^{\perp}$ and $\overline{T^*(H)} = \ker(T)^{\perp}$.

Proof. Both equalities in (a) are immediate consequences of (3.4.1). Using Exercise 3.5 with N = T(H), we get that $\overline{T(H)} = (T(H)^{\perp})^{\perp} = \ker(T^*)^{\perp}$. The second equality in (b) is shown similarly (or by replacing T with T^* in the first one).

Corollary 3.4.13. Let H be a Hilbert space (over \mathbb{F}) and let $T \in \mathcal{B}(H)$. Then T(H) is dense in H if and only if T^* is 1-1.

Proof. Using Proposition 3.4.12 and Corollary 3.2.7, we get

$$\overline{T(H)} = H \iff \ker(T^*)^{\perp} = H \iff \ker(T^*) = \{0\}.$$

As T^* is linear, we also have $\ker(T^*) = \{0\} \Leftrightarrow T^*$ is 1-1.

As another illustration, we finally mention:

Proposition 3.4.14. Let H be a Hilbert space (over \mathbb{F}) and let $T \in \mathcal{B}(H)$. Then T is invertible in $\mathcal{B}(H)$ if and only if T^* is invertible in $\mathcal{B}(H)$, in which case we have $(T^*)^{-1} = (T^{-1})^*$.

Proof. Left to the reader as Exercise 3.19.

3.5 Self-adjoint operators

In this section, we introduce a very important class of bounded operators on a Hilbert space and discuss some of their properties.

Definition 3.5.1. Let H be a Hilbert space (over \mathbb{F}). An operator $T \in \mathcal{B}(H)$ is called *self-adjoint* when $T^* = T$, that is, we have

 $\langle T(x), y \rangle = \langle x, T(y) \rangle$ for all $x, y \in H$.

If $\mathbb{F} = \mathbb{C}$, a self-adjoint operator in $\mathcal{B}(H)$ is also called *Hermitian*, while it is often called *symmetric* if $\mathbb{F} = \mathbb{R}$.

We note that if $T, T' \in \mathcal{B}(H)$ are self-adjoint, and $\lambda \in \mathbb{R}$, then it is obvious that $\lambda T + T'$ is also self-adjoint.

Example 3.5.2. Let $A = [a_{j,k}] \in M_n(\mathbb{F})$ and let $T_A \in \mathcal{B}(\mathbb{F}^n)$ denote the operator given by multiplication with A (cf. Example 3.4.4). Then T_A is self-adjoint if and only if $A^* = A$, i.e., $\overline{a_{k,j}} = a_{j,k}$ for all $j, k \in \{1, \ldots, n\}$. In particular, when $\mathbb{F} = \mathbb{R}$, T_A is self-adjoint if and only if A is symmetric. \Box

Example 3.5.3. Assume *H* is a Hilbert space with a countably infinite orthonormal basis $\mathcal{B} = \{u_j\}_{j \in \mathbb{N}}$.

Let D denote the diagonal operator associated to a bounded sequence $\{\lambda_j\}_{j\in\mathbb{N}}$ in \mathbb{F} (w.r.t. \mathcal{B}), so $D(u_j) = \lambda_j u_j$ for every j, cf. Example 3.4.7. Then D is self-adjoint if and only if $\overline{\lambda_j} = \lambda_j$ for all $j \in \mathbb{N}$, i.e., $\lambda_j \in \mathbb{R}$ for all $j \in \mathbb{N}$. In particular, D is always self-adjoint when $\mathbb{F} = \mathbb{R}$.

Let S denote the right shift operator on H (w.r.t. \mathcal{B}), so $S(u_j) = u_{j+1}$ for every j. As we have seen in Example 3.4.8, S^* is the left shift operator, and it is obvious that $S^* \neq S$. So S is not self-adjoint.

Example 3.5.4. Let (X, \mathcal{A}, μ) be a measure space and set $H := L^2(X, \mathcal{A}, \mu)$. If $f \in \mathcal{L}^{\infty}(X, \mathcal{A}, \mu)$, then the multiplication operator $M_f \in \mathcal{B}(H)$ given by $M_f([g]) = [fg]$, cf. Example 3.4.9, is self-adjoint if and only if $M_{\overline{f}} = M_f$.

Thus, M_f is self-adjoint whenever f is real-valued (μ -a.e.). It can be shown that the converse statement holds whenever (X, \mathcal{A}, μ) satisfies the mild assumption that it is semifinite (cf. Exercise 3.24).

Example 3.5.5. Let $K : [a, b] \times [a, b] \to \mathbb{C}$ be a continuous function and let T_K denote the associated integral operator on $H = L^2([a, b])$, cf. Example 3.4.10, which is determined on C([a, b]) by $(T_K(g))(s) = \int_a^b K(s, t) f(t) dt$. Then T_K is self-adjoint if and only if $T_{K^*} = T_K$ (where $K^*(s, t) = \overline{K(t, s)}$). Hence it is clear that T_K is self-adjoint whenever K is real-valued. We leave it as an exercise to check that the converse statement also holds.

Example 3.5.6. Let M be a closed subspace of a Hilbert space H and let P_M denote the ortogonal projection of H on M. Then P_M is self-adjoint.

Indeed, let $x, y \in H$. As $P_M(x) \in M$ and $y - P_M(y) \in M^{\perp}$, we have

$$\langle P_M(x), y - P_M(y) \rangle = 0.$$

Similarly, we have $\langle x - P_M(x), P_M(y) \rangle = 0$. Hence we get

$$\langle P_M(x), y \rangle = \langle P_M(x), P_M(y) + y - P_M(y) \rangle$$

$$= \langle P_M(x), P_M(y) \rangle + \langle P_M(x), y - P_M(y) \rangle$$

$$= \langle P_M(x), P_M(y) \rangle + \langle x - P_M(x), P_M(y) \rangle$$

$$= \langle P_M(x) + x - P_M(x), P_M(y) \rangle$$

$$= \langle x, P_M(y) \rangle$$

An orthogonal projection map P_M has another property, namely that it is a *positive* operator on H. Such operators are generalizations of operators associated with positive semidefinite matrices. We refer to Exercise 3.29 for the definition and some of the properties of positive operators.

Next, we note that there is an abundance of self-adjoint operators.

Proposition 3.5.7. *Let* H *be a Hilbert space* (over \mathbb{F}) *and* $T \in \mathcal{B}(H)$.

Then the operators $T + T^*$, T^*T and TT^* are all self-adjoint. Moreover, if $\mathbb{F} = \mathbb{C}$, then $\frac{1}{i}(T - T^*)$ is also self-adjoint.

Proof. The reader should have no difficulty to verify these assertions by using the properties of the *-operation on $\mathcal{B}(H)$ listed in Theorem 3.4.2.

A noteworthy consequence is that bounded self-adjoint operators on a *complex* Hilbert space have a canonical decomposition similar to the one enjoyed by complex numbers.

Corollary 3.5.8. Let H be a Hilbert space over \mathbb{C} and let $T \in \mathcal{B}(H)$. Set

$$\operatorname{Re}(T) := \frac{1}{2} (T + T^*), \quad \operatorname{Im}(T) := \frac{1}{2i} (T - T^*).$$

Then $\operatorname{Re}(T)$ and $\operatorname{Im}(T)$ are both self-adjoint, and we have

 $T = \operatorname{Re}(T) + i \operatorname{Im}(T).$

Proof. The first assertion follows readily from Proposition 3.5.7. The second one is elementary.

Consider a bounded operator T on a Hilbert space $H \neq \{0\}$. The *numerical range* of T is defined as the subset of \mathbb{F} given by

$$W_T := \{ \langle T(x), x \rangle : x \in H, ||x|| = 1 \}.$$

Some properties of T are reflected in the geometric properties of W_T and $\overline{W_T}$, cf. Exercise 3.34. We will mainly be interested in the *numerical radius* of T, which is defined by

$$N_T := \sup\{ |\lambda| : \lambda \in W_T \} = \sup\{ |\langle T(x), x \rangle| : x \in H, ||x|| = 1 \}.$$

We note that the Cauchy-Schwarz inequality implies that $N_T \leq ||T||$.

A remarkable fact, proven below, is that N_T agrees with ||T|| when T is self-adjoint. We first observe that if T is self-adjoint, then $W_T \subseteq \mathbb{R}$. Indeed, if $T^* = T$, then for every $x \in H$, we have

$$\langle T(x), x \rangle = \langle x, T(x) \rangle = \langle T(x), x \rangle,$$

and the claim clearly follows.

Theorem 3.5.9. Let H be a Hilbert space, $H \neq \{0\}$, and let $T \in \mathcal{B}(H)$ be self-adjoint. Then we have

$$||T|| = N_T.$$

Proof. It suffices to prove that $||T|| \leq N_T$, hence that

$$||T(x)|| \le N_T \quad \text{for all } x \in H_1. \tag{3.5.1}$$

We first note that if $v \in H$, then $|\langle T(v), v \rangle| \leq N_T ||v||^2$.

Indeed, if v = 0, the claim is trivial. Otherwise, if $v \neq 0$ and $u := \frac{1}{\|v\|} v$, so $v = \|v\| u$, then

$$\left| \langle T(v), v \rangle \right| = \|v\|^2 \left| \langle T(u), u \rangle \right| \leq N_T \|v\|^2.$$

Let now $x \in H_1$. If T(x) = 0, then the inequality in (3.5.1) is trivially satisfied, so we can assume that $T(x) \neq 0$ and set $y := \frac{1}{\|T(x)\|} T(x) \in H_1$. Then we have

$$||T(x)|| = \frac{1}{||T(x)||} \langle T(x), T(x) \rangle = \langle T(x), y \rangle.$$
(3.5.2)

Similarly, $||T(x)|| = \langle y, T(x) \rangle$. As T is self-adjoint, we get

$$||T(x)|| = \langle T(y), x \rangle. \tag{3.5.3}$$

Combining (3.5.2) and (3.5.3), and using our previous observations, as well as the parallellogram law and the fact that $||x|| \leq 1$, ||y|| = 1, we get

$$\|T(x)\| = \frac{1}{2} \left(\langle T(x), y \rangle + \langle T(y), x \rangle \right)$$

$$= \frac{1}{4} \left(\langle T(x+y), x+y \rangle - \langle T(x-y), x-y \rangle \right)$$

$$\leq \frac{1}{4} N_T \left(\|x+y\|^2 + \|x-y\|^2 \right)$$

$$= \frac{1}{2} N_T \left(\|x\|^2 + \|y\|^2 \right)$$

$$\leq N_T.$$

This shows that (3.5.1) is satisfied, as desired.

Having in mind the spectral theorem for symmetric real matrices, it is legitimate to wonder whether it could be true that every self-adjoint operator $T \in \mathcal{B}(H)$ is *diagonalizable* in the sense that there always exists an orthonormal basis for H consisting of eigenvectors for T. However, as the next example illustrates, a self-adjoint operator may not have any eigenvalue at all, so this can not be true in general.

Example 3.5.10. Let $H = L^2([0,1])$ (with usual Lebesgue measure) and let $T = M_f$ be the self-adjoint operator in $\mathcal{B}(H)$ given by multiplication with the bounded continuous function f(t) = t on [0,1], cf. Example 3.5.4. Then we leave it as an exercise (cf. Exercise 3.30) to verify that T has no (complex) eigenvalues.

We will see in the next chapter that every *compact* self-adjoint operator can be diagonalized in the sense mentioned above. Theorem 3.5.9 will help us to make the first step in this direction, by showing that a compact self-adjoint operator T has at least one an eigenvalue, namely ||T|| or -||T||.

3.6 Unitary operators

In this section, we look at another important class of operators on Hilbert spaces. As a warm-up, we first characterize the linear operators which are isometric. If H is a Hilbert space, a map $T: H \to H$ is said to preserve the inner product if it satisfies that

$$\langle T(x), T(y) \rangle = \langle x, y \rangle$$
 for all $x, y \in H$.

Proposition 3.6.1. Let $H \neq \{0\}$ be a Hilbert space (over \mathbb{F}) and let $S: H \to H$. Then the following conditions are equivalent:

- (i) $S \in \mathcal{B}(H)$ and $S^*S = I_H$;
- (*ii*) S is linear and preserves the inner product;
- (*iii*) S is a linear isometry.

Proof. $(i) \Rightarrow (ii)$: Assume $S \in \mathcal{B}(H)$ satisfies $S^*S = I_H$. Then S is linear and for all $x, y \in H$, we have

$$\langle S(x), S(y) \rangle = \langle x, (S^*S)(y) \rangle = \langle x, y \rangle,$$

so (ii) holds.

 $(ii) \Rightarrow (iii)$: Any map preserving the inner product is isometric, so this is evident.

 $(iii) \Rightarrow (i)$: Assume S is a linear isometry. Then $S \in \mathcal{B}(H)$ and $T := S^*S - I_H \in \mathcal{B}(H)$ is self-adjoint. Then for any $x \in H$, we have

$$\langle T(x), x \rangle = \langle (S^*S - I)(x), x \rangle = \langle S(x), S(x) \rangle - \langle x, x \rangle = \|S(x)\|^2 - \|x\|^2 = 0$$

Thus, $W_T = \{0\}$, so, using Theorem 3.5.9, we get that $||T|| = N_T = 0$. Hence, T = 0, i.e., $S^*S = I_H$, so (*i*) holds.

Example 3.6.2. Assume H is finite-dimensional and $S : H \to H$ is a linear isometry, so $S^*S = I_H$, cf. Proposition 3.6.1. As S is injective, it is also surjective (since $\dim(S(H)) = \dim(H) - \dim(\ker(S)) = \dim(H)$, so S(H) = H). Thus, S is bijective, so it has an inverse S^{-1} (which is also a linear isometry). Since $S^*S = I_H$, we get that $S^{-1} = S^*$. In particular, we also have $SS^* = I_H$.

Remark 3.6.3. When *H* is infinite-dimensional, then a linear isometry *S* is not necessarily surjective. A typical example is the right shift operator *S* considered in Example 3.4.6, whose range does not contain the first basis vector; in this case, we have $S^*S = I_H$, while $SS^* \neq I_H$ (cf. Exercise 3.22).

Definition 3.6.4. Let H be a Hilbert space (over \mathbb{F}). An operator $U \in \mathcal{B}(H)$ is called *unitary* when it satisfies

$$U^*U = UU^* = I_H.$$

Thus, $U \in \mathcal{B}(H)$ is unitary if and only if U is bijective and $U^{-1} = U^*$.

When $\mathbb{F} = \mathbb{R}$, it is customary to say *orthogonal* instead of unitary.

Proposition 3.6.5. Let H be a Hilbert space (over \mathbb{F}) and let $U : H \to H$. Then the following conditions are equivalent:

(i) $U \in \mathcal{B}(H)$ and U is unitary;

- (ii) U is bijective, linear and preserves the inner product;
- (iii) U is a surjective linear isometry.

Proof. $(i) \Rightarrow (ii)$: If $U \in \mathcal{B}(H)$ is unitary, then U is bijective and linear, and Proposition 3.6.1 gives that it preserves the inner product. Hence, (ii) holds.

 $(ii) \Rightarrow (iii)$: This implication is evident.

 $(iii) \Rightarrow (i)$: Suppose U is a surjective linear isometry. As a linear isometry is injective, U is bijective. Moreover, Proposition 3.6.1 gives that $U^*U = I_H$. So we get that $U^{-1} = U^*$, i.e., U is unitary, and (i) holds.

Example 3.6.6. Assume *H* has a countably infinite orthonormal basis \mathcal{B} and *D* is the diagonal operator associated to a bounded sequence $\{\lambda_j\}_{j\in\mathbb{N}}$ in \mathbb{F} (w.r.t. \mathcal{B}).

Then it is straightforward to check that D is unitary if and only if $\overline{\lambda_j}\lambda_j = 1$, i.e., $|\lambda_j| = 1$, for all $j \in \mathbb{N}$.

Example 3.6.7. Let (X, \mathcal{A}, μ) be a measure space and set $H := L^2(X, \mathcal{A}, \mu)$. For $f \in \mathcal{L}^{\infty}$, consider the multiplication operator $M_f \in \mathcal{B}(H)$. Then we clearly have

$$(M_f)^* M_f = M_{|f|^2} = M_f (M_f)^*,$$

so we see that M_f is unitary whenever |f| = 1 μ -a.e. The converse holds if μ is semifinite, cf. Exercise 3.31.

Example 3.6.8. Let $H = \ell^2(\mathbb{Z})$. We may then define the bilateral forward shift operator $U : H \to H$ by

$$[U(\xi)](j) = \xi(j-1)$$
 for all $\xi \in H$ and all $j \in \mathbb{Z}$.

Indeed, since $\sum_{j\in\mathbb{Z}} |\xi(j-1)|^2 = \sum_{j\in\mathbb{Z}} |\xi(j)|^2 < \infty$, we see that $U(\xi) \in H$ and $||U(\xi)||_2 = ||\xi||_2$ for every $\xi \in H$. Thus U is isometric.

We may now conclude from Proposition 3.6.5 that U is unitary. Its adjoint $U^* = U^{-1}$ is called *the bilateral backward shift operator* (on $H = \ell^2(\mathbb{Z})$). We note that if $\mathcal{B} = \{e_n\}_{n \in \mathbb{Z}}$ denotes the canonical basis of $H = \ell^2(\mathbb{Z})$ as in Example 3.3.4, then we have

$$U(e_n) = e_{n+1}$$
 and $U^*(e_n) = e_{n-1}$ for all $n \in \mathbb{Z}$.

Example 3.6.9. Let $H = L^2(\mathbb{R}, \mathcal{A}, \mu)$, where \mathcal{A} denote the Lebesguemeasurable subsets of \mathbb{R} and μ is the usual Lebesgue measure on \mathcal{A} . For each $x \in \mathbb{R}$ and $f \in \mathcal{L}^2(\mathbb{R}, \mathcal{A}, \mu)$, define $f_x : \mathbb{R} \to \mathbb{C}$ by $f_x(y) = f(y - x)$ for every $y \in \mathbb{R}$. Then f_x is Lebesgue-measurable, and, using the translation invariance of μ , we get that

$$||f_x||_2^2 = \int_{\mathbb{R}} |f_x(y)|^2 d\mu(y) = \int_{\mathbb{R}} |f(y-x)|^2 d\mu(y) = \int_{\mathbb{R}} |f(y')|^2 d\mu(y') = ||f||_2^2.$$

It follows that the map $U_x : H \to H$ given by $U_x([f]) = [f_x]$ for each f in $L^2(\mathbb{R}, \mathcal{A}, \mu)$ is a well-defined unitary operator on H, satisfying that $(U_x)^* = U_{-x}$ (and $U_{x+x'} = U_x U_{x'}$ for all $x, x' \in \mathbb{R}$). You are asked to provide all the missing details in Exercise 3.33.

Let now H, K be Hilbert spaces (over \mathbb{F}). A bijective, linear map U from H onto K which preserves the inner product is often called an *isomorphism* of Hilbert spaces. As in Proposition 3.6.5, one shows that it is equivalent to require that U is a surjective linear isometry, or that $U \in \mathcal{B}(H, K)$ is unitary in the sense that we have $U^*U = I_H$ and $UU^* = I_K$. (Here, $U^* \in \mathcal{B}(K, H)$ denotes the adjoint of U, cf. Remark 3.4.3). We will therefore say that H and K are *isomorphic as Hilbert spaces* when such a map $U : H \to K$ exists.

Theorem 3.6.10. Let $H \neq \{0\}$ be a Hilbert space over \mathbb{C} , and let \mathcal{B} be an orthonormal basis of H. Then H and $\ell^2(\mathcal{B})$ are isomorphic as Hilbert spaces.

Proof. Let $x \in H$ and define $\hat{x} : \mathcal{B} \to \mathbb{C}$ by

$$\widehat{x}(u) := \langle x, u \rangle \quad \text{for all } u \in \mathcal{B}.$$

Note that Parseval's identity says that $\sum_{u \in \mathcal{B}} |\hat{x}(u)|^2 = ||x||^2$. In particular, we have $\hat{x} \in \ell^2(\mathcal{B})$ and $||\hat{x}|| = ||x||$. Thus we can define an isometric map $U: H \to \ell^2(\mathcal{B})$ by

$$U(x) = \hat{x}$$
 for all $x \in H$.

It is elementary to check that U is linear. Moreover, U is surjective. Indeed, let $\xi \in \ell^2(\mathcal{B})$. As $\sum_{u \in \mathcal{B}} |\xi(u)|^2 < \infty$, the set

$$\mathcal{B}_{\xi} := \{ u \in \mathcal{B} : \xi(u) \neq 0 \}$$

must be countable. Let $\{u_j\}_{j\in N}$ be an enumeration of \mathcal{B}_{ξ} , where $N = \{1, \ldots, n\}$ for some $n \in \mathbb{N}$ or $N = \mathbb{N}$. Note that if $N = \mathbb{N}$, we have $\sum_{j=1}^{\infty} |\xi(u_j)|^2 < \infty$, and this implies readily that the sequence $\left\{\sum_{j=1}^{k} \xi(u_j) u_j\right\}_{k\in\mathbb{N}}$ is Cauchy, hence convergent in H.

Thus we may define $x \in H$ by $x := \sum_{j \in N} \xi(u_j) u_j$, and we then have

$$\widehat{x}(u) = \langle x, u \rangle = \sum_{j \in N} \xi(u_j) \langle u_j, u \rangle = \begin{cases} \xi(u_k) & \text{if } u = u_k \text{ for some } k \in N, \\ 0 & \text{if } u \in \mathcal{B} \setminus \mathcal{B}_{\xi}, \end{cases}$$

i.e., $\hat{x}(u) = \xi(u)$ for all $u \in \mathcal{B}$. Hence, $U(x) = \xi$, showing that U is surjective.

We can now conclude that U is an isomorphism of Hilbert spaces from H to $\ell^2(\mathcal{B})$, as we wanted to show.

Remark 3.6.11. Theorem 3.6.10 is also true when $H \neq \{0\}$ is a Hilbert space over \mathbb{R} , but one has then to replace $\ell^2(\mathcal{B})$ with the real ℓ^2 -space

$$\ell^{2}_{\mathbb{R}}(\mathcal{B}) := \left\{ \xi : \mathcal{B} \to \mathbb{R} : \sum_{u \in \mathcal{B}} |\xi(u)|^{2} < \infty, \right\}$$

considered as a Hilbert space over \mathbb{R} .

Remark 3.6.12. If $H \neq \{0\}$ is a Hilbert space over \mathbb{C} , and $\mathcal{B}, \mathcal{B}'$ are both orthonormal bases of H, then we get from Theorem 3.6.10 that $\ell^2(\mathcal{B})$ and $\ell^2(\mathcal{B}')$ are isomorphic as Hilbert spaces. It can be shown that this implies that (and in fact is equivalent to) \mathcal{B} and \mathcal{B}' having the same cardinality, meaning that there exists a bijection between \mathcal{B} and \mathcal{B}' . (A similar statement holds if $H \neq \{0\}$ is a Hilbert space over \mathbb{R}).

3.7 Exercises

In the exercises of this chapter, H always denotes a Hilbert space over \mathbb{F} , unless otherwise stated.

Exercise 3.1. Let X be an inner product space. Check that the parallelogram law and the polarization identities hold.

Exercise 3.2. Let H be a Hilbert space which is infinite-dimensional, i.e., is not finite-dimensional (as a vector space). Argue first that there exists an orthonormal sequence $\{x_n\}_{n\in\mathbb{N}}$ in H. Then use this sequence to show that the unit ball H_1 is not compact.

Exercise 3.3. Consider $X := \ell^{\infty}(\mathbb{N})$ as a metric space w.r.t. $d(f,g) = ||f - g||_u$. Let A be the subset of X given by

$$A := \left\{ a^{(N)} : N \in \mathbb{N} \right\},\$$

where $a^{(N)}(n) = 1$ if $1 \le n \le N$ and $a^{(N)}(n) = 0$ if n > N.

a) Show that A is closed in X.

b) Let $x \in X$ be given by x(n) = 1 + 1/n for all $n \in \mathbb{N}$.

Show that d(x, A) = 1 and that $1 < d(x, a^{(N)})$ for all $N \in \mathbb{N}$.

Exercise 3.4. Let $c \in H$, r > 0 and set $B := B_r(c) = \{y \in H : ||y-c|| \le r\}$. Check that B is closed and convex, and give a formula for x_B when $x \in H \setminus B$.

Exercise 3.5. Let S denote a nonempty subset of H.

a) Show that S^{\perp} is a closed subspace of H.

b) Set $M := \overline{\text{Span}(S)}$. Verify that $S^{\perp} = M^{\perp}$. Then deduce that $M = (S^{\perp})^{\perp}$. Deduce also that if N is a subspace of H, then $\overline{N} = (N^{\perp})^{\perp}$.

Exercise 3.6. Let M be a closed subspace of H and $x \in H$.

- a) Check that the associated map $P_M: H \to H$ is linear.
- b) Show that $P_M(x) = y$ for some $y \in M$ if and only if $x y \in M^{\perp}$.

Then show that $P_M(x)$ is the unique vector y in M such that $x - y \in M^{\perp}$.

Exercise 3.7. Assume $P \in \mathcal{B}(H)$ satisfies that $P^2 = P$ and ||P|| = 1.

Show that P(H) is closed and $H = P(H) \oplus \ker P$. Then show that P is the orthogonal projection of H on M := P(H).

Exercise 3.8. Let (X, \mathcal{A}, μ) be a measure space.

a) Show that

 $\langle [f], [g] \rangle := \int_X f \,\overline{g} \, \mathrm{d}\mu$

gives a well-defined inner product on $L^2 := L^2(X, \mathcal{A}, \mu)$ (cf. Example 3.1.12).

b) Let $E \in \mathcal{A}$. Set $\mathcal{A}_E = \{A \cap E : A \in \mathcal{A}\}$ and $\mu_E = \mu_{|\mathcal{A}_E}$. We recall that $(E, \mathcal{A}_E, \mu_E)$ is a measure space.

Show that there exists an isometric isomorphism from $L^2(E, \mathcal{A}_E, \mu_E)$ onto the space M_E defined in Example 3.2.9, i.e.,

 $M_E = \left\{ [g] : g \in \mathcal{L}^2 \text{ and } g \text{ lives essentially on } E \right\}.$

Exercise 3.9. Consider $H = L^2([a, b], \mathcal{A}, \mu)$, where \mathcal{A} denotes the σ -algebra of all Lebesgue measurable subsets of [a, b], and μ is the usual Lebesgue measure on \mathcal{A} . Set

$$M := \left\{ [g] \in H : g \in \mathcal{L}^2([a,b],\mathcal{A},\mu), \int_{[a,b]} g \, \mathrm{d}\mu = 0 \right\}$$

Check that M is a closed subspace of H. Then, given $[f] \in H$, find an expression for the best approximation of [f] in M.

Exercise 3.10. Let $H \neq \{0\}$. Show that the following conditions are equivalent:

- (a) H is separable;
- (b) There is a sequence satisfying the assumptions in Example 3.3.5;
- (c) H has a countable orthonormal basis.

Note that Example 3.3.5 shows that (b) \Rightarrow (c). So it suffices to show that (a) \Rightarrow (b), and (c) \Rightarrow (a).

Exercise 3.11. In the context of Fourier analysis described in Example 3.3.12 (see also Example 3.3.2), the formula

$$||f||_{2}^{2} = \sum_{n \in \mathbb{Z}} \left| \widehat{f}(n) \right|^{2}$$

is called *Parseval's identity*. (The more general equality obtained in Theorem 3.3.8 c) is also often called Parseval's identity.)

a) Set f(t) = t for all $t \in [-\pi, \pi]$. Compute the Fourier coefficients of f.

b) Use a) and Parseval's identity to show that $\sum_{n=1}^{\infty} \frac{1}{n^2} = \frac{\pi^2}{6}$.

c) Set $g(t) = e^t$ for all $t \in [-\pi, \pi]$. Use Parseval's identity to obtain a formula for the sum of the series $\sum_{n=1}^{\infty} \frac{1}{n^2+1}$.

Exercise 3.12. Let $H = L^2([-1, 1], \mathcal{A}, \mu)$, where \mathcal{A} denote the Lebesguemeasurable subsets of [-1, 1] and μ is the restriction of the usual Lebesgue measure to \mathcal{A} .

For each $n \in \{0\} \cup \mathbb{N}$, let $p_{n+1} : [-1,1] \to \mathbb{C}$ be defined by $p_{n+1}(t) = t^n$, and set $S := \{[p_{n+1}] : n \in \{0\} \cup \mathbb{N}\} \subseteq H$.

a) Show that Span(S) is dense in H.

b) Apply the Gram-Schmidt orthonormalization process to S to obtain an orthonormal basis $\mathcal{B} = \{[q_{n+1}] : n \in \{0\} \cup \mathbb{N}\}$ for H, where each q_{n+1} is the polynomial on [-1, 1] given by

$$q_{n+1}(t) = \frac{\sqrt{n+\frac{1}{2}}}{2^n n!} \frac{d^n}{dt^n} \left((t^2 - 1)^n \right).$$

(These polynomials are called the normalized Legendre polynomials.)

Exercise 3.13. Let H be the L^2 -space on $[-\pi, \pi]$ w.r.t. to the normalized Lebesgue measure μ , as in Example 3.3.2. Set

$$H_{\mathsf{even}} := \{ [f] \in H : f \text{ is even} \} \text{ and } H_{\mathsf{odd}} := \{ [f] \in H : f \text{ is odd} \}.$$

We recall that a function $f: [-\pi, \pi] \to \mathbb{C}$ is called *even* if f(-t) = f(t) for all t, while it is called *odd* if f(-t) = -f(t) for all t.

a) Show that H_{even} is a closed subspace of H and that $(H_{\text{even}})^{\perp} = H_{\text{odd}}$. Then describe the orthogonal projection P of H on H_{even} .

Hint: It might be helpful to consider the map $[f] \to [\tilde{f}]$, where $\tilde{f}(t) := f(-t)$.

b) Find an orthonormal basis for H_{even} and one for H_{odd} .

Exercise 3.14. Let H_1, H_2 be Hilbert spaces over \mathbb{F} and consider

$$H := H_1 \times H_2 = \{ (x_1, x_2) : x_1 \in H_1, x_2 \in H_2 \}$$

as a vector space over \mathbb{F} when equipped with its natural pointwise operations. For $(x_1, x_2), (y_1, y_2) \in H$, set

$$\langle (x_1, x_2), (y_1, y_2) \rangle := \langle x_1, y_1 \rangle + \langle x_2, y_2 \rangle.$$

a) Check that this gives an inner product on H such that H is a Hilbert space. Check also that the associated norm on H is given by

$$||(x_1, x_2)|| = (||x_1||^2 + ||x_2||^2)^{1/2}.$$
b) Set $\widetilde{H}_1 := \{(x_1, 0) : x_1 \in H_1\}$ and $\widetilde{H}_2 := \{(0, x_2) : x_2 \in H_2\},\$

Check that \widetilde{H}_1 and \widetilde{H}_2 are closed subspaces of H and that \widetilde{H}_j is isomorphic to H_j (j = 1, 2).

Check also that $H = \widetilde{H}_1 \oplus \widetilde{H}_2$ and describe the orthogonal projection from H on \widetilde{H}_1 (resp. \widetilde{H}_2).

Note: The Hilbert space H is called the external direct sum of H_1 and H_2 .

Exercise 3.15. Let H_1 and H_2 be Hilbert spaces over \mathbb{F} , and let H be the external direct sum of H_1 and H_2 , as defined in Exercise 3.14. Assume \mathcal{B}_1 and \mathcal{B}_2 are orthonormal bases for H_1 and H_2 , respectively.

Find an orthonormal basis \mathcal{B} for H in terms of \mathcal{B}_1 and \mathcal{B}_2 .

Exercise 3.16. The concept of generalized sums can be used to provide an alternative way of describing Fourier expansions in Hilbert spaces.

Let X be a normed space, J be a nonempty set, $\{x_j\}_{j\in J}$ be a family of elements of X, and $x \in X$. Then one says that the generalized sum $\sum_{j\in J} x_j$ converges to x when the following holds: given $\varepsilon > 0$, there exists a finite subset $F_0 \subseteq J$ such that for all finite subsets F of J containing F_0 , we have

$$\left\| x - \sum_{j \in F} x_j \right\| < \varepsilon \,,$$

in which case we write

$$x = \sum_{j \in J} x_j \, .$$

Consider a Hilbert space H and $x \in H$.

a) Show that we have

$$x = \sum_{u \in \mathcal{B}} \langle x, u \rangle \, u \, .$$

b) Show also that if M is a closed subspace of H and C is an orthonormal basis for M, then we have

$$P_M(x) = \sum_{v \in \mathcal{C}} \langle x, v \rangle v.$$

Exercise 3.17. Let $T \in \mathcal{B}(H)$. Assume H_0 is a dense subspace of H which is invariant under T, and let $T_0 \in \mathcal{B}(H_0)$ denote the restriction of T to H_0 . Further, assume there exists some $S_0 \in \mathcal{B}(H_0)$ such that

$$\langle T_0(x), y \rangle = \langle x, S_0(y) \rangle$$
 for all $x, y \in H_0$.

Show that $T^* = S$, where $S \in \mathcal{B}(H)$ is the unique extension of S_0 provided by Theorem 1.4.2.

Exercise 3.18. Show that the formula for $(T_K)^*$ in Example 3.4.10 is correct.

Hint : Consider $H_0 = [g] : g \in C([a, b])$ and use Exercise 3.17.

Exercise 3.19. Prove Proposition 3.4.14.

Exercise 3.20. Let $v, w \in H$ and consider the linear operator $T_{v,w} : H \to H$ defined by

$$T_{v,w}(x) := \langle x, v \rangle w$$
 for all $x \in H$.

Note that $T_{v,w}$ has rank one if $v, w \in H \setminus \{0\}$.

a) Show that $T_{v,w}$ is bounded with norm $||T_{v,w}|| = ||v|| ||w||$. Then show that $(T_{v,w})^* = T_{w,v}$.

b) Show that every $T \in \mathcal{B}(H)$ which has rank one is of the form $T = T_{v,w}$ for some $v, w \in H \setminus \{0\}$.

c) Assume $T \in \mathcal{B}(H)$ is a finite-rank operator, $T \neq 0$. Show that T may be written as a finite sum of rank one operators in $\mathcal{B}(H)$.

Hint : Start by picking an orthonormal basis for T(H).

d) Show that if $T \in \mathcal{B}(H)$ is a finite-rank operator, then so is T^* .

Exercise 3.21. Let $T \in \mathcal{B}(H)$ and let M be a closed subspace of H. Show that

M is invariant under T if and only if M^{\perp} is invariant under T^* .

Exercise 3.22. Let $T \in \mathcal{B}(H)$.

a) Show that $\ker(T) = \ker(T^*T)$ and $\overline{T^*(H)} = \overline{(T^*T)(H)}$.

b) Assume T is normal, i.e., satisfies $T^*T = TT^*$. Show that

$$\ker(T^*) = \ker(T)$$
 and $T^*(H) = T(H)$.

c) Assume T is normal and has an eigenvalue λ . Show that $\overline{\lambda}$ is an eigenvalue of T^* , and that $E_{\overline{\lambda}}^{T^*} = E_{\lambda}^T$.

d) Assume *H* has a countably infinite orthonormal basis $\mathcal{B} = \{u_j\}_{j \in \mathbb{N}}$ and let $S \in \mathcal{B}(H)$ be the right shift operator (w.r.t. \mathcal{B}). Set $T = S^*$.

Check that $TT^* = S^*S = I_H$, while $T^*T = SS^* = P$, where P is the orthogonal projection of H on $\{u_1\}^{\perp}$. Deduce that T is not normal.

Next, check that 0 is an eigenvalue for T, while 0 is not an eigenvalue of $T^* = S$. (*Note*: This shows that the assertion in c) does not necessarily hold when T is not normal.)

Finally, if you are in the right mood, show that S has no eigenvalues, while every λ satisfying $|\lambda| < 1$ is an eigenvalue of T.

Exercise 3.23. Let *H* and *K* be Hilbert spaces over \mathbb{F} , and let $T \in \mathcal{B}(H, K)$.

a) Show that there exists a unique operator $T^* \in \mathcal{B}(K, H)$ (called the adjoint of T) satisfying that

$$\langle T(x), y \rangle = \langle x, T^*(y) \rangle$$
 for all $x \in H$ and all $y \in K$.

b) Let $T' \in \mathcal{B}(H, K)$ and $\alpha, \beta \in \mathbb{F}$. Let also L be a Hilbert space over \mathbb{F} , and let $S \in \mathcal{B}(K, L)$, so that $ST \in \mathcal{B}(H, L)$. Show that the following properties hold:

• i) $(\alpha T + \beta T')^* = \overline{\alpha} T^* + \overline{\beta} T'^*;$ ii) $(ST)^* = T^*S^*;$ iii) $(T^*)^* = T;$ • iv) $||T^*|| = ||T||;$ v) $||T^*T|| = ||T||^2.$

Exercise 3.24. Let (X, \mathcal{A}, μ) be a measure space. One says that (X, \mathcal{A}, μ) is *semifinite* when the following condition holds: if $E \in \mathcal{A}$ and $\mu(E) = \infty$, then there exists $F \subseteq E$, $F \in \mathcal{A}$ such that $0 < \mu(F) < \infty$.

a) Show that (X, \mathcal{A}, μ) is semifinite whenever it is σ -finite.

Assume from now on that (X, \mathcal{A}, μ) is semifinite. Set $H := L^2(X, \mathcal{A}, \mu)$. Let $f \in \mathcal{L}^{\infty}$ and consider the multiplication operator $M_f \in \mathcal{B}(H)$ defined in Example 3.5.4.

b) Show that $||M_f|| = ||f||_{\infty}$.

c) Show that if M_f is self-adjoint, then f is real-valued μ -a.e. (As observed in Example 3.4.9, the converse is true without any restriction on (X, \mathcal{A}, μ) .)

Exercise 3.25. Assume $P \in \mathcal{B}(H)$ is a self-adjoint projection, i.e., it satisfies that $P^* = P = P^2$. Show that P is the orthogonal projection of H on M := P(H) (which is a closed subspace of H).

Exercise 3.26. Let $H \neq \{0\}$.

a) Assume $T \in \mathcal{B}(H)$ is self-adjoint. Deduce from Theorem 3.5.9 that T = 0 if and only if $\langle T(x), x \rangle = 0$ for all $x \in H$.

b) Suppose $\mathbb{F} = \mathbb{R}$. Give an example with $H = \mathbb{R}^2$ showing that the equivalence in a) may fail when T is not self-adjoint.

c) Assume $\mathbb{F} = \mathbb{C}$ and let $T \in \mathcal{B}(H)$. Show that T = 0 if and only if $\langle T(x), x \rangle = 0$ for all $x \in H$.

Exercise 3.27. Show that the set $\mathcal{B}(H)_{sa} := \{T \in \mathcal{B}(H) : T^* = T\}$ is closed in $\mathcal{B}(H)$.

Exercise 3.28. Let $H \neq \{0\}$. If $T \in \mathcal{B}(H)$ is self-adjoint, we have seen that $W_T \subseteq \mathbb{R}$; of course, if $\mathbb{F} = \mathbb{R}$, this gives no information on T as this inclusion is then true for any T in $\mathcal{B}(H)$. We assume therefore in this exercise that $H \neq \{0\}$ is a Hilbert space over \mathbb{C} .

Let $T \in \mathcal{B}(H)$. Then show that the following assertions are equivalent:

- (i) T is self-adjoint;
- (*ii*) $W_T \subseteq \mathbb{R}$;
- (*iii*) $\langle T(x), x \rangle \in \mathbb{R}$ for all $x \in H$.

Exercise 3.29. A self-adjoint operator T in $\mathcal{B}(H)$ is called *positive* when

$$\langle T(x), x \rangle \ge 0 \quad \text{for all } x \in H,$$

$$(3.7.1)$$

in which case we write $T \ge 0$.

(We note that if $\mathbb{F} = \mathbb{C}$ and $T \in \mathcal{B}(H)$ satisfies (3.7.1), then T is automatically self-adjoint, as follows from the previous exercise.)

a) Let $S \in \mathcal{B}(H)$, and let $R \in \mathcal{B}(H)$ be self-adjoint.

Check that $S^*S \ge 0$ and $R^2 \ge 0$. Then show that

$$||S|| \le 1 \Leftrightarrow (I_H - S^*S) \ge 0.$$

- b) Let M be a closed subspace of H. Check that $P_M \ge 0$.
- c) Assume that $T, T' \in \mathcal{B}(H)$ are positive and $\lambda \in [0, \infty)$.

Check that T + T' and λT are also positive.

d) Show that the set of positive operators in $\mathcal{B}(H)$ is closed in $\mathcal{B}(H)$.

Exercise 3.30. Let $H = L^2([0, 1])$ (with usual Lebesgue measure) and let $T = M_f$ be the self-adjoint operator in $\mathcal{B}(H)$ given by multiplication with the function f(t) = t on [0, 1], cf. Example 3.5.4. Show that T has no (complex) eigenvalues.

Exercise 3.31. Let (X, \mathcal{A}, μ) be a semifinite measure space (cf. Exercise 3.24), and let $f \in \mathcal{L}^{\infty}$. Suppose that the multiplication operator M_f on $H = L^2(X, \mathcal{A}, \mu)$ is unitary. Then show that $|f| = 1 \mu$ -a.e.

(As observed in Example 3.6.7, the converse statement is true without any restriction on (X, \mathcal{A}, μ) .)

Exercise 3.32. Assume $H \neq \{0\}$ is separable (cf. Exercise 3.10) and infinitedimensional. Let then \mathcal{B} be an orthonormal basis for H indexed by \mathbb{Z} , say $\mathcal{B} = \{v_k\}_{k \in \mathbb{Z}}$. One may then define the bilateral shift operator $V : H \to H$ (w.r.t. \mathcal{B}) by

$$V(x) = \sum_{k \in \mathbb{Z}} \langle x, v_k \rangle v_{k+1} \text{ for all } x \in H, \text{ i.e., by}$$
$$V(x) = \lim_{n \to \infty} \sum_{k=-n}^{n} \langle x, v_k \rangle v_{k+1} \text{ for all } x \in H.$$

a) Show that V is a unitary operator in $\mathcal{B}(H)$.

b) Describe V as a multiplication operator when $H = L^2([-\pi, \pi])$ (with normalized Lebesgue measure μ) and $v_k(t) = e^{ikt}$ for every $k \in \mathbb{Z}$.

c) Assume $\mathbb{F} = \mathbb{C}$. Let $U : H \to \ell^2(\mathbb{Z})$ denote the isomorphism of Hilbert spaces defined in the proof of Theorem 3.6.10. Show that UVU^* is the bilateral forward shift operator on $\ell^2(\mathbb{Z})$.

Exercise 3.33. Provide all the details missing in Example 3.6.9.

Exercise 3.34. Let T be a bounded operator on a Hilbert space $H \neq \{0\}$. Check that the following properties of W_T and N_T hold:

- (a) $W_{T^*} = \left\{ \overline{\lambda} : \lambda \in W_T \right\}$; hence, $N_{T^*} = N_T$.
- (b) W_T contains all the possible eigenvalues of T.
- (c) $W_{\alpha T+\beta I_H} = \alpha W_T + \beta$ for all $\alpha, \beta \in \mathbb{F}$.
- (d) $W_{UTU^*} = W_T$, hence $N_{UTU^*} = N_T$, for every unitary $U \in \mathcal{B}(H)$.
- (e) W_T is a compact subset of \mathbb{F} when H is finite-dimensional.

It can also be shown that W_T is a convex subset of \mathbb{F} . This result is called the Toeplitz-Hausdorff Theorem, but the proof is beyond the scope of these notes.

CHAPTER 4

On compact operators

4.1 Introduction to compact operators between normed spaces

In this section, X and Y will denote normed spaces, both over \mathbb{F} , unless otherwise specified.

Definition 4.1.1. An operator $T \in \mathcal{L}(X, Y)$ is called *compact* if the sequence $\{T(x_n)\}_{n \in \mathbb{N}}$ has a convergent subsequence in Y whenever $\{x_n\}_{n \in \mathbb{N}}$ is a bounded sequence in X.

We set $\mathcal{K}(X,Y) := \{T \in \mathcal{L}(X,Y) : T \text{ is compact } \}.$

To appreciate this definition, the concept of relative compactness for subsets of a metric space will be helpful.

A subset A of a metric space (Z, d) is called *relatively compact in* Z if its closure \overline{A} is a compact subset of Z. (Some authors say *precompact* instead of relatively compact.) Equivalently, and this may be taken as the definition for our purposes, a subset A of Z is relatively compact in Z if and only if every sequence in A has a subsequence which converges in Z. In comparison, we recall that A is compact if and only if every sequence in A has a subsequence which converges in A.

We also remark that a subset A of Z is bounded whenever A is relatively compact in Z: indeed, if A is not bounded, then we can pick (any) $z \in Z$ and find a sequence $\{a_n\}_{n\in\mathbb{N}}$ in A such that $d(a_n, z) > n$ for all $n \in \mathbb{N}$; it is then rather easy to see that $\{a_n\}_{n\in\mathbb{N}}$ has no convergent subsequence in Z, so A is not relatively compact.

Proposition 4.1.2. Let $T \in \mathcal{L}(X, Y)$. Then T is compact if and only if T(B) is relatively compact in Y whenever B is a bounded subset of X.

Proof. Assume first that T is compact and let $B \subseteq X$ be bounded. We want to show that T(B) is relatively compact in Y. So let $\{y_n\}_{n\in\mathbb{N}}$ be a sequence in T(B). For each $n \in \mathbb{N}$ we may then write $y_n = T(x_n)$ for some $x_n \in B$. As the sequence $\{x_n\}_{n\in\mathbb{N}}$ lies in B, it is bounded. Hence, by compactness of T, $\{y_n\}_{n\in\mathbb{N}} = \{T(x_n\}_{n\in\mathbb{N}} \text{ has a convergence subsequence in } Y$. Thus, T(B)is relatively compact, as desired.

Conversely, assume that T maps bounded subsets of X into relatively compact subsets of Y. We want to show that T is compact. So let $\{x_n\}_{n \in \mathbb{N}}$ be a bounded sequence in X. Set $B := \{x_n : n \in \mathbb{N}\}$. Then B is a bounded subset of X, so $T(B) = \{T(x_n) : n \in \mathbb{N}\}$ is relatively compact in Y. As $\{T(x_n)\}_{n \in \mathbb{N}}$ is a sequence in T(B), we can conclude that it has a convergent subsequence in Y. Thus, T is compact, as desired.

Corollary 4.1.3. Assume $T \in \mathcal{L}(X, Y)$ is compact. Then T is bounded. Thus, $\mathcal{K}(X, Y) \subseteq \mathcal{B}(X, Y)$.

Proof. Set $B := X_1$. Since B is a bounded subset of X, we get from Proposition 4.1.2 that T(B) is relatively compact subset of Y. This implies that T(B) is bounded. Hence we can find M > 0 such that $||T(x)|| \le M$ for all $x \in X_1$, and it follows that T is bounded with $||T|| \le M$.

An important class of compact operators consists of the finite-rank operators in $\mathcal{B}(X, Y)$. We recall that $T \in \mathcal{L}(X, Y)$ is said to have finite-rank if T(X) is finite-dimensional.

Proposition 4.1.4. Assume that $T \in \mathcal{B}(X, Y)$ has finite-rank. Then T is compact.

Proof. Assume $\{x_n\}_{n\in\mathbb{N}}$ is a bounded sequence in X, and let M > 0 be such that $||x_n|| \leq M$ for all $n \in \mathbb{N}$. Then we have

$$||T(x_n)|| \le ||T|| ||x_n|| \le ||T|| M$$

for all $n \in \mathbb{N}$. Now, the ball $B := \{y \in Y : ||y|| \le ||T|| M\}$ is closed in Y. Considering T(X) as a normed space w.r.t. to the norm it inherits from Y, we get that the set $K := T(X) \cap B$ is a closed and bounded subset of T(X). Since T(X) is finite-dimensional (by assumption), we get from Proposition 1.3.2 that K is compact in T(X). As $\{T(x_n)\}_{n\in\mathbb{N}}$ is a sequence in K, we can therefore conclude that it has a convergent subsequence. This shows that T is compact. **Example 4.1.5.** Consider $X = C([0, 1], \mathbb{R})$ with the uniform norm $\|\cdot\|_u$. For $g \in X$, define $T(g) : [0, 1] \to \mathbb{R}$ by

$$[T(g)](s) = \int_0^1 \sin(s-t) g(t) \, \mathrm{d}t \quad \text{for all } s \in [0,1].$$

Since $\sin(s-t) = \sin(s)\cos(t) - \cos(s)\sin(t)$, we have that

$$[T(g)](s) = \left(\int_0^1 \cos(t) g(t) \, dt\right) \,\sin(s) - \left(\int_0^1 \sin(t) g(t) \, dt\right) \,\cos(s)$$

for all $s \in [0, 1]$. It follows that $T(g) \in X$. Moreover, the map $T : X \to X$ sending g to T(g) is clearly linear. As T(X) is 2-dimensional, T has finite-rank. Further, since

$$\left| [T(g)](s) \right| \le \int_0^1 |\sin(s-t)g(t)| \, \mathrm{d}t \le \int_0^1 |g(t)| \, \mathrm{d}t \le ||g||_u$$

for all $s \in [0, 1]$, we get that $||T(g)||_u \leq ||g||_u$ for all $g \in X$. Hence, T is bounded. We can therefore conclude that T is compact.

More generally, using the Arzelà-Ascoli Theorem (cf. Lindstrøm's book), it can be shown that if a function $K : [a, b] \times [c, d] \to \mathbb{R}$ is continuous, then the associated integral operator $T : C([c, d], \mathbb{R}) \to C([a, b], \mathbb{R})$, defined by

$$[T(g)](s) = \int_{c}^{d} K(s,t) g(t) \, \mathrm{d}t \quad \text{for all } s \in [a,b] \,,$$

is compact.

Theorem 4.1.6. $\mathcal{K}(X,Y)$ is a subspace of $\mathcal{B}(X,Y)$. Moreover, if Y is a Banach space, then $\mathcal{K}(X,Y)$ is closed in $\mathcal{B}(X,Y)$, and it follows that $\mathcal{K}(X,Y)$ is a Banach space.

Proof. We leave the proof of the first assertion as an exercise. So assume that Y is Banach space, and let $\{T_n\}_{n\in\mathbb{N}}$ be a sequence in $\mathcal{K}(X,Y)$ converging to some $T \in \mathcal{B}(X,Y)$. To show that $\mathcal{K}(X,Y)$ is closed in $\mathcal{B}(X,Y)$, we have to show that T is compact.

So let $\{x_n\}_{n\in\mathbb{N}}$ be a bounded sequence in X. Choose M > 0 such that $||x_n|| \leq M$ for all $n \in \mathbb{N}$.

• Since T_1 is compact, there exists a subsequence $\{x_{n_k}\}_{k\in\mathbb{N}}$ of $\{x_n\}_{n\in\mathbb{N}}$ such that $T_1(x_{n_k}) \to y_1$ as $k \to \infty$ for some $y_1 \in Y$.

We set $x_{1,k} := x_{n_k}$ for each $k \in \mathbb{N}$. We then have $T_1(x_{1,n}) \to y_1$ as $n \to \infty$.

- Similarly, since $\{x_{1,n}\}_{n\in\mathbb{N}}$ is bounded and T_2 is compact, we can find a sequence $\{x_{2,n}\}_{n\in\mathbb{N}}$, which is a subsequence of $\{x_{1,n}\}_{n\in\mathbb{N}}$, and therefore of $\{x_n\}_{n\in\mathbb{N}}$, such that $T_2(x_{2,n}) \to y_2$ as $n \to \infty$ for some $y_2 \in Y$.
- Proceeding inductively, for each $m \in \mathbb{N}$, $m \geq 2$, we can find a sequence $\{x_{m,n}\}_{n\in\mathbb{N}}$, which is a subsequence of $\{x_{m-1,n}\}_{n\in\mathbb{N}}$, and therefore of $\{x_n\}_{n\in\mathbb{N}}$, such that $T_m(x_{m,n}) \to y_m$ as $n \to \infty$ for some $y_m \in Y$.

We now set $x'_k := x_{k,k} \in X$ for each $k \in \mathbb{N}$, and claim that

$${T(x'_k)}_{k\in\mathbb{N}}$$
 is a Cauchy sequence in Y. (4.1.1)

Since Y is complete, we will then be able to conclude that $\{T(x'_k)\}_{k\in\mathbb{N}}$ is convergent, hence that $\{T(x_n)\}_{n\in\mathbb{N}}$, which will show that T is compact.

To establish (4.1.1), we first observe that for any $k, l, m \in \mathbb{N}$, we have

$$\begin{aligned} \|T(x_{l}') - T(x_{k}')\| &\leq \|(T - T_{m})(x_{l}') + T_{m}(x_{l}') - T_{m}(x_{k}') + (T_{m} - T)(x_{k}')\| \\ &\leq \|(T - T_{m})(x_{l}')\| + \|T_{m}(x_{l}') - T_{m}(x_{k}')\| + \|(T_{m} - T)(x_{k}')\| \\ &\leq \|T - T_{m}\| \|x_{l}'\| + \|T_{m}(x_{l}') - T_{m}(x_{k}')\| + \|T_{m} - T\| \|x_{k}'\| \\ &\leq \|T_{m}(x_{l}') - T_{m}(x_{k}')\| + 2M \|T - T_{m}\|. \end{aligned}$$

Let then $\varepsilon > 0$ and choose $m \in \mathbb{N}$ such that $||T - T_m|| < \varepsilon/3M$. By construction, for each $k \ge m$, we have that $T_m(x'_k) = T_m(x_{k,k})$ is an element of the sequence $\{T_m(x_{m,n})\}_{n\in\mathbb{N}}$, which is convergent to y_m . It follows that the sequence $\{T_m(x'_k)\}_{k\in\mathbb{N}}$ is convergent, hence that it is Cauchy. So we can pick $N \in \mathbb{N}$ such that $||T_m(x'_l) - T_m(x'_k)|| < \varepsilon/3$ for all $k, l \ge N$. This gives that

$$\|T(x_{l}') - T(x_{k}')\| \le \|T_{m}(x_{l}') - T_{m}(x_{k}')\| + 2M \|T - T_{m}\| < \varepsilon/3 + 2M (\varepsilon/3M) = \varepsilon$$

for all $k, l \ge N$. Hence we have shown that the claim (4.1.1) is true.

Finally, as Y is a Banach space, we know that $\mathcal{B}(X, Y)$ is a Banach space too, and this implies that $\mathcal{K}(X, Y)$, being closed in $\mathcal{B}(X, Y)$, is also a Banach space.

An immediate consequence is the following:

Corollary 4.1.7. Assume Y is a Banach space and set

$$\mathcal{F}(X,Y) := \{T \in \mathcal{B}(X,Y) : T \text{ has finite-rank}\}.$$

Then we have

$$\overline{\mathcal{F}(X,Y)} \subseteq \mathcal{K}(X,Y).$$

Example 4.1.8. Let $1 \leq p < \infty$ and set $X := \ell^p(\mathbb{N})$, which we know is a Banach space w.r.t. $\|\cdot\|_p$. For each $\lambda \in \ell^\infty(\mathbb{N})$, we may consider the multiplication operator $M_\lambda \in \mathcal{B}(X)$ given by

$$[M_{\lambda}(x)](n) = \lambda(n) x(n)$$

for all $x \in X$ and all $n \in \mathbb{N}$. One readily checks that $||M_{\lambda}|| = ||\lambda||_{\infty}$.

Now, assume that $\lambda \in c_0(\mathbb{N})$, i.e., $\lim_{n\to\infty} \lambda(n) = 0$. Then M_{λ} is compact. Indeed, for each $k \in \mathbb{N}$, let $\lambda^{(k)} \in \ell^{\infty}(\mathbb{N})$ be defined by

$$\lambda^{(k)}(n) = \begin{cases} \lambda(n) & \text{if } 1 \le k \le n, \\ 0 & \text{otherwise,} \end{cases}$$

for every $n \in \mathbb{N}$. Then it is clear that each $M_{\lambda^{(k)}}$ has finite-rank; moreover,

$$||M_{\lambda} - M_{\lambda^{(k)}}|| = ||\lambda - \lambda^{(k)}||_{\infty} \to 0 \quad \text{as } k \to \infty.$$

Thus $M_{\lambda} \in \overline{\mathcal{F}(X, X)} \subseteq \mathcal{K}(X, X)$.

We set $\mathcal{K}(X) := \mathcal{K}(X, X)$, so that $\mathcal{K}(X)$ is a subspace of $\mathcal{B}(X)$. If X is finite-dimensional, then every operator in $\mathcal{B}(X)$ has finite-rank, so $\mathcal{K}(X) = \mathcal{B}(X)$. On the other hand, if X is infinite-dimensional, then $\mathcal{K}(X) \neq \mathcal{B}(X)$, the reason being that the identity operator I_X is not compact in this case: indeed, if X is infinite-dimensional, then $I_X(X_1) = X_1$ is closed, but not compact, (cf. Exercise 1.2).

We also mention (cf. Exercise 4.1) that $\mathcal{K}(X)$ is a two-sided ideal in $\mathcal{B}(X)$, meaning that we have

$$ST \in \mathcal{K}(X)$$
 if $S \in \mathcal{B}(X)$ and $T \in \mathcal{K}(X)$, or if $S \in \mathcal{K}(X)$ and $T \in \mathcal{B}(X)$.

This property implies that no operator in $\mathcal{K}(X)$ can have a bounded inverse when X is infinite-dimensional (because if some $T \in \mathcal{K}(X)$ had an inverse $T^{-1} \in \mathcal{B}(X)$, then we would have that $I_X = T^{-1}T \in \mathcal{K}(X)$, hence that $\dim(X) < \infty$).

We end this section with an interesting result concerning the possible eigenvalues of a compact operator, and their associated eigenspaces.

Theorem 4.1.9. Let $T \in \mathcal{K}(X)$. Then the following facts hold:

(a) Let $\delta > 0$. Then $\{\lambda \in \mathbb{F} : \lambda \text{ is an eigenvalue of } T \text{ and } |\lambda| > \delta\}$ is a finite set.

(b) If $\lambda \in \mathbb{F}$ is a non-zero eigenvalue of T, then the associated eigenspace $E_{\lambda} := \{x \in X : T(x) = \lambda x\}$ is finite-dimensional.

(c) The set of eigenvalues of T (which may be empty) is countable and bounded. If this set is countably infinite and $\{\lambda_k : k \in \mathbb{N}\}$ is an enumeration of it, then $\lim_{k\to\infty} \lambda_k = 0$.

As we will be mostly interested in compact *self-adjoint* operators acting on Hilbert spaces in this course, for which much more can be said (cf. Theorem 4.3.4), we skip the proof of this theorem.

4.2 Compact operators on Hilbert spaces

In view of Corollary 4.1.7, it is natural to wonder whether any compact operator from a normed space to a Banach space may be approximated in operator norm by bounded finite-rank operators. This problem was open until 1973, when a counterexample was exhibited by P. Enflo. Happily, the situation is as nice as possible when the target space is a Hilbert space.

Theorem 4.2.1. Let X be a normed space and H be a Hilbert space (both over \mathbb{F}). Then we have

$$\overline{\mathcal{F}(X,H)} = \mathcal{K}(X,H) \,.$$

Proof. By Corollary 4.1.7, we only have to show that $\mathcal{K}(X, H) \subseteq \overline{\mathcal{F}(X, H)}$. So let $T \in \mathcal{K}(X, H)$ and $\varepsilon > 0$. We need to prove that there exists $S \in \mathcal{F}(X, H)$ such that $||T - S|| \leq \varepsilon$. Clearly, we can assume $T \neq 0$.

Set $A := \overline{T(X_1)}$. Since X_1 is bounded and T is compact, the set A is compact in H. As H is a metric space, this implies that A is totally bounded (cf. Proposition 3.5.12 in Lindstrøm's book). Hence we can cover A with some open balls B_1, \ldots, B_n of radius $\varepsilon/4$, having respective centers $a_1, \ldots, a_n \in A$. For each $j = 1, \ldots, n$, we can then find $x_j \in X_1$ such that $||a_j - T(x_j)|| < \varepsilon/4$.

Set now $F := \text{Span}(\{T(x_1), \ldots, T(x_n)\})$, which is a finite dimensional subspace of H, and let P_F denote the orthogonal projection of H on F. Since the range of P_FT is contained in F, P_FT has finite-rank, so $P_FT \in \mathcal{F}(X, H)$. We claim that

$$||T - P_F T|| \le \varepsilon.$$

Indeed, let $x \in X_1$. Then $T(x) \in A$, so $T(x) \in B_j$ for some $j \in \{1, \ldots, n\}$. Hence,

$$||T(x) - T(x_j)|| \le ||T(x) - a_j|| + ||a_j - T(x_j)|| < \varepsilon/4 + \varepsilon/4 = \varepsilon/2.$$

Further, since $T(x_j) \in F$, we have that $P_F(T(x_j)) = T(x_j)$. Thus, using also that $||P_F|| = 1$, we obtain that

$$\|(T - P_F T)(x)\| = \|T(x) - T(x_j) + (P_F T)(x_j) - (P_F T)(x)\|$$

$$\leq \|T(x) - T(x_j)\| + \|P_F(T(x_j) - T(x))\|$$

$$\leq \|T(x) - T(x_j)\| + \|P_F\|\|T(x_j) - T(x)\|$$

$$= 2 \|T(x) - T(x_j)\|$$

$$< 2 \cdot \varepsilon/2 = \varepsilon.$$

As this holds for every $x \in X_1$, the claim follows. Hence, setting $S := P_F T$, we are done.

Remark 4.2.2. Let X be a normed space and H be a Hilbert space, and let $T \in \mathcal{K}(X, H)$. Then it can be shown that $\overline{T(X)}$ is separable. We leave this as an exercise.

Theorem 4.2.1 immediately gives:

Corollary 4.2.3. Let H be a Hilbert space. Set $\mathcal{K}(H) := \mathcal{K}(H, H)$ and $\mathcal{F}(H) := \mathcal{F}(H, H)$. Then we have

$$\overline{\mathcal{F}(H)} = \mathcal{K}(H) \,.$$

An application of this result is the following:

Corollary 4.2.4. Let H be a Hilbert space and let $T \in \mathcal{K}(H)$. Then $T^* \in \mathcal{K}(H)$. In other words, $\mathcal{K}(H)$ is closed under the adjoint operation.

Proof. Using Corollary 4.2.3, we can find a sequence $\{T_n\}_{n\in\mathbb{N}}$ in $\mathcal{F}(H)$ such that $||T - T_n|| \to 0$ as $n \to \infty$. Now, it is not difficult to see that $\mathcal{F}(H)$ is closed under the adjoint operation (cf. Exercise 3.20). Hence, $\{T_n^*\}_{n\in\mathbb{N}}$ is a sequence in $\mathcal{F}(H)$, and we have

$$||T^* - T_n^*|| = ||(T - T_n)^*|| = ||T - T_n|| \to 0 \text{ as } n \to \infty.$$

Thus, $T^* \in \overline{\mathcal{F}(H)} = \mathcal{K}(H)$.

We recall from the previous section that if H is finite-dimensional, then $\mathcal{K}(H) = \mathcal{B}(H) = \mathcal{F}(H)$, while $\mathcal{K}(H) \neq \mathcal{B}(H)$ if H is infinite-dimensional. An elementary argument showing that I_H is not compact when H is infinitedimensional goes as follows: letting $\{u_j\}_{j\in\mathbb{N}}$ be any orthonormal sequence in H, we have $||u_j - u_k|| = \sqrt{2}$ for all $j, k \in \mathbb{N}$, and it follows that the sequence $\{I_H(u_j)\}_{j\in\mathbb{N}} = \{u_j\}_{j\in\mathbb{N}}$ does not have any convergent subsequence.

On compact operators

Let H be an infinite-dimensional Hilbert space H. An interesting class of compact operators on H containing $\mathcal{F}(H)$ consists of the so-called Hilbert-Schmidt operators. For simplicity, we only consider the case where H is separable. We note that every orthonormal basis for H is then countable: indeed, assume (for contradiction) that H had an uncountable orthonormal basis \mathcal{B} . Then, as $||u - u'|| = \sqrt{2}$ for all distinct $u, u' \in \mathcal{B}$, we see that any dense subset of H would have to be uncountable, contradicting the separability of H.

Lemma 4.2.5. Assume H is a separable, infinite-dimensional Hilbert space (over \mathbb{F}). Let $\mathcal{B} = \{u_j\}_{j \in \mathbb{N}}$ and $\mathcal{C} = \{v_j\}_{j \in \mathbb{N}}$ be orthonormal bases for H, and let $T \in \mathcal{B}(H)$. Then we have

$$\sum_{j=1}^{\infty} \|T(u_j)\|^2 = \sum_{j=1}^{\infty} \|T(v_j)\|^2.$$

Proof. Using Parseval's identity (two times), we get

$$\sum_{j=1}^{\infty} ||T(u_j)||^2 = \sum_{j=1}^{\infty} \sum_{k=1}^{\infty} \left| \langle T(u_j), v_k \rangle \right|^2 = \sum_{j=1}^{\infty} \sum_{k=1}^{\infty} \left| \langle u_j, T^*(v_k) \rangle \right|^2$$
$$= \sum_{j=1}^{\infty} \sum_{k=1}^{\infty} \left| \langle T^*(v_k), u_j \rangle \right|^2 = \sum_{k=1}^{\infty} \sum_{j=1}^{\infty} \left| \langle T^*(v_k), u_j \rangle \right|^2$$
$$= \sum_{k=1}^{\infty} ||T^*(v_k)||^2.$$

Note that the change of order of summation at the second but last step above is allowed since we are dealing with sums of non-negative numbers. Applying what we have done to the case where $\mathcal{B} = \mathcal{C}$, i.e., $u_j = v_j$ for every $j \in \mathbb{N}$, we get that

$$\sum_{j=1}^{\infty} \|T(v_j)\|^2 = \sum_{k=1}^{\infty} \|T^*(v_k)\|^2.$$

Thus we obtain that

$$\sum_{j=1}^{\infty} \|T(u_j)\|^2 = \sum_{k=1}^{\infty} \|T^*(v_k)\|^2 = \sum_{j=1}^{\infty} \|T(v_j)\|^2,$$

as desired.

Remark 4.2.6. An analogous result is true when H is finite-dimensional and \mathcal{B}, \mathcal{C} are orthonormal bases for H.

Definition 4.2.7. Let H be a separable, infinite-dimensional Hilbert space (over \mathbb{F}). An operator $T \in \mathcal{B}(H)$ is called an *Hilbert-Schmidt operator* when we have

$$\sum_{j=1}^{\infty} \|T(u_j)\|^2 < \infty$$

for some orthonormal basis $\mathcal{B} = \{u_i\}_{i \in \mathbb{N}}$ of H, in which case we set

$$||T||_2 := \left(\sum_{j=1}^{\infty} ||T(u_j)||^2\right)^{1/2}.$$

Lemma 4.2.5 shows that the definition of T being a Hilbert-Schmidt operator, and the value of $||T||_2$, do not depend on the choice of orthonormal basis for H. We set

 $\mathcal{HS}(H) := \{T \in \mathcal{B}(H) : T \text{ is a Hilbert-Schmidt operator}\}.$

Proposition 4.2.8. Let H be a separable, infinite-dimensional Hilbert space (over \mathbb{F}).

Then $\mathcal{HS}(H)$ is a subspace of $\mathcal{K}(H)$, which contains $\mathcal{F}(H)$ and is closed under the adjoint operation.

Moreover, the map $T \to ||T||_2$ is a norm on $\mathcal{HS}(H)$, which satisfies

$$||T|| \le ||T||_2$$

for every $T \in \mathcal{HS}(H)$.

Proof. We first note that it is evident from the proof of Lemma 4.2.5 that $T^* \in \mathcal{HS}(H)$ whenever $T \in \mathcal{HS}(H)$.

Let $\mathcal{B} = \{u_j\}_{j \in \mathbb{N}}$ be an orthonormal basis for H, and let $T, T' \in \mathcal{HS}(H)$. Define $\xi, \xi' \in \ell^2(\mathbb{N})$ by

$$\xi(j) := ||T(u_j)||$$
 and $\xi'(j) := ||T'(u_j)||$ for each $j \in \mathbb{N}$,

so that $\|\xi\|_2 = \|T\|_2$ and $\|\xi'\|_2 = \|T'\|_2$. Using the triangle inequality, first in H, and then in $\ell^2(\mathbb{N})$, we get

$$\sum_{j=1}^{\infty} \|(T+T')(u_j)\|^2 \le \sum_{j=1}^{\infty} \left(\|T(u_j)\| + \|T'(u_j)\| \right)^2 = \|\xi + \xi'\|_2^2$$
$$\le (\|\xi\|_2 + \|\xi'\|_2)^2 = (\|T\|_2 + \|T'\|_2)^2 < \infty.$$

This shows that $T + T' \in \mathcal{HS}(H)$ and

$$||T + T'||_2 \le ||T||_2 + ||T'||_2.$$

Moreover, one easily checks that $\lambda T \in \mathcal{HS}(H)$ and $\|\lambda T\|_2 = |\lambda| \|T\|_2$ for every $\lambda \in \mathbb{F}$. If $\|T\|_2 = 0$, then we get that $\|T(u_j)\| = 0$ for every $j \in \mathbb{N}$, and this clearly implies that T = 0.

Hence, we have shown so far that $\mathcal{HS}(H)$ is a subspace of $\mathcal{B}(H)$ which is closed under the adjoint operation, and that $\|\cdot\|_2$ is a norm on $\mathcal{HS}(H)$.

To show that $||T|| \leq ||T||_2$, let $x \in H \setminus \{0\}$. Set $v_1 = \frac{1}{||x||}x$ and let $\{v_j\}_{j\geq 2}$ be an orthonormal basis for $\{x\}^{\perp}$. Then $\{v_j\}_{j\in\mathbb{N}}$ is an orthonormal basis for H, so we get

$$\|T(x)\|^{2} = \|x\|^{2} \|T(v_{1})\|^{2} \le \|x\|^{2} \sum_{j=1}^{\infty} \|T(v_{j})\|^{2} = \|T\|_{2}^{2} \|x\|^{2}.$$

us. $\|T\| \le \|T\|_{2}.$

Thus, $||T|| \le ||T||_2$.

Next, we show that $T \in \mathcal{K}(H)$. For each $n \in \mathbb{N}$, let P_n denote the orthogonal projection of H on Span $(\{u_1, \ldots u_n\})$ and set $T_n := TP_n$. Then we have

$$\sum_{j=1}^{\infty} \|T_n(u_j)\|^2 = \sum_{j=1}^n \|T(u_j)\|^2 < \infty,$$

so $T_n \in \mathcal{HS}(H)$ for each $n \in \mathbb{N}$. Hence,

$$||T - T_n|| \le ||T - T_n||_2 = \left(\sum_{j=n+1}^{\infty} ||T(u_j)||^2\right)^{1/2} \to 0 \text{ as } n \to \infty.$$

Since $T_n \in \mathcal{F}(H)$ for each *n*, Theorem 4.1.7 gives that $T \in \mathcal{K}(H)$. Hence, $\mathcal{HS}(H) \subseteq \mathcal{K}(H)$.

It only remains to show that $\mathcal{F}(H) \subseteq \mathcal{HS}(H)$, but we leave this as an exercise.

Remark 4.2.9. For additional properties of $\mathcal{HS}(H)$, see Exercise 4.6.

Remark 4.2.10. If $H \neq \{0\}$ is finite-dimensional and $\mathcal{B} = \{u_j\}_{j=1}^n$ is an orthonormal basis for H, then we get a norm on $\mathcal{B}(H)$ by setting

$$||T||_2 := \left(\sum_{j=1}^n ||T(u_j)||^2\right)^{1/2}$$

(which does not depend on the choice of orthonormal basis for H).

Letting $A = [a_{i,j}]$ denotes the matrix of T w.r.t. \mathcal{B} , one readily checks that

$$||T||_2 = \left(\sum_{i,j=1}^n |a_{i,j}|^2\right)^{1/2},$$

i.e., $||T||_2$ coincides with the so-called Fröbenius-norm of A.

Example 4.2.11. Set $H = L^2([a, b], \mathcal{A}, \mu)$, where \mathcal{A} denotes the Lebesgue measurable subsets of a closed interval [a, b] and μ is the Lebesgue measure on \mathcal{A} . Let $K : [a, b] \times [a, b] \to \mathbb{C}$ be a continuous function and let $T_K \in \mathcal{B}(H)$ denote the associated integral operator on H, which is the extension of the integral operator $T_K : C([a, b]) \to C([a, b])$ given by

$$[T_K(f)](s) = \int_a^b K(s,t)f(t) \, \mathrm{d}t \quad \text{for } f \in C([a,b]) \text{ and } s \in [a,b].$$

cf. Example 2.1.9 and Exercise 2.10. Then T_K is a Hilbert-Schmidt operator on H (so T_K is compact by Proposition 4.2.8).

To show this, we start by picking an orthonormal basis $\mathcal{B} = \{[u_j]\}_{j \in \mathbb{N}}$ for H, where each u_j is a continuous functions on [a, b]. (One may for example construct \mathcal{B} by applying the Gram-Schmidt orthonormalization process to the monomials $\{t^{j-1}: j \in \mathbb{N}\}$). We note that $\overline{\mathcal{B}} := \{[\overline{u_j}]\}_{j \in \mathbb{N}}$ is then also an orthonormal basis for H.

Let now $s \in [a, b]$ and let $k_s \in C([a, b])$ be given by $k_s(t) = K(s, t)$ for all $t \in [a, b]$. Note that for each $j \in \mathbb{N}$, we have

$$[T_K(u_j)](s) = \int_a^b K(s,t) \, u_j(t) \, \mathrm{d}t = \int_{[a,b]} k_s(t) \, \overline{\overline{u_j}(t)} \, \mathrm{d}\mu(t) = \langle [k_s], [\overline{u_j}] \rangle \,.$$

Moreover, Parseval's identity gives that

$$\| [k_s] \|_2 = \left(\sum_{j=1}^{\infty} \left| \langle [k_s], [\overline{u_j}] \rangle \right|^2 \right)^{1/2}.$$

Thus, we obtain that

$$\sum_{j=1}^{\infty} \left| [T_K(u_j)](s) \right|^2 = \sum_{j=1}^{\infty} \left| \langle [k_s], [\overline{u_j}] \rangle \right|^2 = \| [k_s] \|_2^2$$

Now, using this and the Monotone Convergence Theorem, we get

$$\begin{split} \sum_{j=1}^{\infty} \| T_K([u_j]) \|_2^2 &= \sum_{j=1}^{\infty} \int_{[a,b]} \left| [T_K(u_j)](s) \right|^2 \, \mathrm{d}\mu(s) \\ &= \int_{[a,b]} \left(\sum_{j=1}^{\infty} \left| [T_K(u_j)](s) \right|^2 \right) \, \mathrm{d}\mu(s) \\ &= \int_{[a,b]} \| [k_s] \|_2^2 \, \mathrm{d}\mu(s) \\ &= \int_{[a,b]} \left(\int_{[a,b]} |k_s(t)|^2 \, \mathrm{d}\mu(t) \right) \, \mathrm{d}\mu(s) \\ &= \int_a^b \int_a^b |K(s,t)|^2 \, \mathrm{d}t \, \mathrm{d}s < \infty, \end{split}$$

which shows that $T_K \in \mathcal{HS}(H)$ with $||T_K||_2 \leq (\int_a^b \int_a^b |K(s,t)|^2 \, \mathrm{d}s \, \mathrm{d}t)^{1/2}$.

In the previous example, one may allow the kernel K to be discontinuous and still obtain an Hilbert-Schmidt operator T_K , as long as K is squareintegrable w.r.t. to the product Lebesgue measure on $[a, b] \times [a, b]$. However, this requires a better knowledge of measure theory than we assume in these notes.

4.3 The spectral theorem for a compact self-adjoint operator

Throughout this section we let H denote a Hilbert space (over \mathbb{F}) different from $\{0\}$. Our main goal is to generalize the spectral theorem for symmetric real matrices known from linear algebra, and prove that every compact self-adjoint compact operator T on H is diagonalizable in the sense that there exists an orthonormal basis for H consisting of eigenvectors of T.

We begin with a series of lemmas.

Lemma 4.3.1. Assume $T \in \mathcal{K}(H)$ has a nonzero eigenvalue $\lambda \in \mathbb{F}$. Then the associated eigenspace $E_{\lambda} := \ker(T - \lambda I)$ is finite-dimensional.

Proof. Assume for contraction that E_{λ} is infinite-dimensional. We may then find a sequence $\{v_n\}_{n\in\mathbb{N}}$ of unit vectors in E_{λ} which are pairwise orthogonal. By compactness of T, $\{T(v_n)\}_{n\in\mathbb{N}}$ has a convergent subsequence. So we may as well assume that $\{T(v_n)\}_{n\in\mathbb{N}}$ is convergent, hence that it is a Cauchy sequence. However, we have that

$$||T(v_n) - T(v_m)||^2 = ||\lambda v_n - \lambda v_m||^2 = 2 |\lambda|^2 \neq 0$$

for all $m, n \in \mathbb{N}$. So $\{T(v_n)\}_{n \in \mathbb{N}}$ is not a Cauchy sequence, giving a contradiction.

Lemma 4.3.2. Let $T \in \mathcal{B}(H)$ be self-adjoint, and assume T has an eigenvalue $\lambda \in \mathbb{F}$. Then $\lambda \in \mathbb{R}$.

Moreover, if λ' is an eigenvalue of T distinct from λ , then $E_{\lambda} \perp E_{\lambda'}$, i.e., $\langle x, y \rangle = 0$ whenever $x \in E_{\lambda}$ and $y \in E_{\lambda'}$.

Proof. Let $x \in E_{\lambda}$. If ||x|| = 1, then we have

$$\lambda = \lambda \langle x, x \rangle = \langle \lambda x, x \rangle = \langle T(x), x \rangle \in W_T \subseteq \mathbb{R},$$

so $\lambda \in \mathbb{R}$. Moreover, assume that λ' is an eigenvalue of T distinct from λ , and let $y \in E_{\lambda'}$. Then we have that $\lambda' \in \mathbb{R}$, so

$$\lambda \langle x, y \rangle = \langle T(x), y \rangle = \langle x, T(y) \rangle = \lambda' \langle x, y \rangle.$$

Since $\lambda' \neq \lambda$, we get that $\langle x, y \rangle = 0$.

Lemma 4.3.3. Let $T \in \mathcal{K}(H)$ be self-adjoint. Then T has an eigenvalue $\lambda \in \mathbb{R}$ such that $|\lambda| = ||T||$.

Proof. If T = 0, then the assertion is trivial. So assume that $T \neq 0$. Using Theorem 3.5.9, we can find a sequence $\{x_n\}_{n\in\mathbb{N}}$ of unit vectors in H such that $|\langle T(x_n), x_n \rangle| \to ||T||$ as $n \to \infty$. Since $\langle T(x_n), x_n \rangle \in \mathbb{R}$ for every n, we can assume (by passing to a subsequence and relabelling) that

$$\langle T(x_n), x_n \rangle \to \lambda \text{ as } n \to \infty, \text{ where } \lambda = \pm ||T||.$$
 (4.3.1)

Moreover, since T is compact, we can also assume (by passing again to a subsequence and relabelling) that $T(x_n) \to y$ as $n \to \infty$ for some $y \in H$. Note that the Cauchy-Schwarz inequality gives that

$$|\langle T(x_n), x_n \rangle| \leq ||T(x_n)||$$
 for every $n \in \mathbb{N}$,

so, letting $n \to \infty$, we get that $||y|| \ge |\lambda| > 0$, so $y \ne 0$.

Now, using that T is self-adjoint, λ is real, $||x_n|| = 1$, and (4.3.1), we get

$$\|T(x_n) - \lambda x_n\|^2 = \langle T(x_n) - \lambda x_n, T(x_n) - \lambda x_n \rangle$$

= $\|T(x_n)\|^2 - 2\lambda \langle T(x_n), x_n \rangle + \lambda^2 \|x_n\|^2$
 $\leq \|T\|^2 - 2\lambda \langle T(x_n), x_n \rangle + \lambda^2$
= $2\lambda \left(\lambda - \langle T(x_n), x_n \rangle\right)$
 $\rightarrow 0 \text{ as } n \rightarrow \infty.$

Thus, $||T(x_n) - \lambda x_n|| \to 0$ as $n \to \infty$, and this gives that

$$||y - \lambda x_n|| \le ||y - T(x_n)|| + ||T(x_n) - \lambda x_n|| \to 0 \text{ as } n \to \infty.$$

Hence,

$$T(y) = \lim_{n \to \infty} T(\lambda x_n) = \lambda \lim_{n \to \infty} T(x_n) = \lambda y$$

Since $y \neq 0$, λ is an eigenvalue of T, as we wanted to show.

We are now ready for the spectral theorem for a compact self-adjoint operator T. Intuitively, we could hope to be able to construct an orthonormal basis of eigenvectors for T by using Lemma 4.3.3 repeatedly as follows:

Start by picking a unit eigenvector v_0 of T associated to the eigenvalue $\lambda_0 = \pm ||T||$. Next, consider the restriction T_1 of T to the closed subspace $M_1 := \{v_0\}^{\perp}$. Pick a unit eigenvector $v_1 \in M_1$ of T_1 associated to the eigenvalue $\lambda_1 = \pm ||T_1||$. Then continue this process inductively.

There are several technicalities involved in working out the details of this approach. We will follow a more pedestrian route, which also provides more information about T.

Theorem 4.3.4. Let $T \in \mathcal{K}(H)$ be self-adjoint, $T \neq 0$. Then there exists an orthonormal basis \mathcal{E} for H which consists of eigenvectors of T.

More precisely, the following facts hold:

- (a) The set L consisting of all nonzero eigenvalues of T is a nonempty, countable subset of the interval $\left[-\|T\|, \|T\|\right]$, containing $\|T\|$ or $-\|T\|$.
- (b) If L is countably infinite, and $\{\lambda_k : k \in \mathbb{N}\}$ is an enumeration of L, then we have $\lim_{k\to\infty} \lambda_k = 0$.
- (c) The eigenspace $E_{\lambda} := \ker(T \lambda I)$ is finite-dimensional for each $\lambda \in L$. (The dimension of E_{λ} is called the (geometric) multiplicity of λ .)
- (d) For each $\lambda \in L$, let \mathcal{E}_{λ} be an orthonormal basis for E_{λ} , and set

$$\mathcal{E}' := \bigcup_{\lambda \in L} \mathcal{E}_{\lambda} \,.$$

Then \mathcal{E}' is an orthonormal basis for $\overline{T(H)} = \ker(T)^{\perp}$, which is countable.

- (e) If ker(T) = {0}, set $\mathcal{E}_0 := \emptyset$; otherwise, let \mathcal{E}_0 be an orthonormal basis for ker(T). Then $\mathcal{E} := \mathcal{E}_0 \cup \mathcal{E}'$ is an orthonormal basis for H which consists of eigenvectors of T.
- (f) Let P_{λ} denote the orthogonal projection of H on E_{λ} for each $\lambda \in L$. Then $P_{\lambda}P_{\lambda'} = 0$ whenever $\lambda \neq \lambda'$ belong to L. Moreover, T has a spectral decomposition

 $T = \sum_{\lambda \in L} \lambda P_{\lambda} \quad (w.r.t. \text{ operator norm}), \text{ meaning that} \\ \lim_{n \to \infty} \|T - \sum_{k=1}^{n} \lambda_k P_{\lambda_k}\| = 0 \quad \text{if } L \text{ is countably infinite} \\ \text{and } \{\lambda_k : k \in \mathbb{N}\} \text{ is an enumeration of } L, \text{ as in } (b).$

Proof. (a): The set L is a subset of \mathbb{R} by Lemma 4.3.2, which contains ||T|| or -||T|| by Lemma 4.3.3. If $\lambda \in L$, and v is an associated eigenvector in H_1 , we have

$$|\lambda| = |\langle \lambda v, v \rangle| = |\langle T(v), v \rangle| \le ||T||$$

Thus, $L \subseteq [-\|T\|, \|T\|]$

To show that L is countable, let $\varepsilon > 0$ and consider the subset of L given by $L_{\varepsilon} := \{\lambda \in L : |\lambda| \ge \varepsilon\}$. Then L_{ε} is finite.

Indeed, assume L_{ε} is nonempty. Then for each $\lambda \in L$, we can pick $v_{\lambda} \in H_1$ such that $T(v_{\lambda}) = \lambda v_{\lambda}$; for $\lambda, \lambda' \in L_{\varepsilon}, \lambda \neq \lambda'$, we then have $\lambda v_{\lambda} \perp \lambda' v_{\lambda'}$ by Lemma 4.3.2, so we get

$$||T(v_{\lambda}) - T(v_{\lambda'})||^2 = ||\lambda v_{\lambda} - \lambda' v_{\lambda'}||^2 = |\lambda|^2 + |\lambda'|^2 \ge 2\varepsilon^2$$

Hence, if L_{ε} was infinite, we could find a sequence in H_1 which T maps into a sequence with no convergent subsequence, contradicting the compactness of T. Thus, L_{ε} is finite.

Now, since $L = \bigcup_{n \in \mathbb{N}} L_{1/n}$, it follows that L is countable.

(b): Assume L is countably infinite and $\{\lambda_k : k \in \mathbb{N}\}$ is an enumeration of L. Let $\varepsilon > 0$ be given. Then, as in (a), we get that the set $K := \{k \in \mathbb{N} : |\lambda_k| \ge \varepsilon\}$ is finite. So there exists $N \in \mathbb{N}$ such that $K \subseteq \{1, \ldots, N\}$. For every $k \ge N + 1$, we then have $|\lambda_k| < \varepsilon$. This shows that $\lim_{k\to\infty} \lambda_k = 0$.

- (c): This is a consequence of Lemma 4.3.1.
- (d): We first remark that since T is self-adjoint, we have

$$\overline{T(H)} = \overline{T^*(H)} = (\ker T)^{\perp}.$$

Next, it follows from Lemma 4.3.2 that $\mathcal{E}_{\lambda} \perp \mathcal{E}_{\lambda'}$ whenever $\lambda \neq \lambda'$ belong to L. So it is clear that \mathcal{E}' is an orthonormal set in H, which is countable since each \mathcal{E}_{λ} is finite and L is countable. Hence, \mathcal{E}' is a countable orthonormal basis for $M := \overline{\text{Span}(\mathcal{E}')}$, and it remains only to show that $M = \ker(T)^{\perp}$, i.e., that $M^{\perp} = \ker(T)$.

• $\underline{\ker(T) \subseteq M^{\perp}}$: Assume $y \in \ker(T)$. Then for each $\lambda \in L$ and $v \in \mathcal{E}_{\lambda}$, we have

 $\lambda \langle v, y \rangle = \langle T(v), y \rangle = \langle v, T(y) \rangle = \langle v, 0 \rangle = 0.$

Since $\lambda \neq 0$, this shows that $y \in (\mathcal{E}')^{\perp} = M^{\perp}$.

• $\underline{M^{\perp} \subseteq \ker(T)}$: It is easy to check that M is invariant under T. Hence, $\overline{M^{\perp}}$ is invariant under $T^* = T$ (cf. Exercise 3.21). We may therefore consider the restriction S of T to M^{\perp} . Then $S \in \mathcal{K}(M^{\perp})$: if not, then there would exist a bounded sequence in M^{\perp} , hence in H, which is mapped by S, hence by T, to a sequence with no convergent subsequence, contradicting the compactness of T. Moreover, S is self-adjoint (this is an easy exercise).

Now, assume that $S \neq 0$. Then Lemma 4.3.3 gives that S has an nonzero eigenvalue μ . This implies that μ is a nonzero eigenvalue of T, i.e., $\mu \in L$. But if $v \in M^{\perp}$ is an eigenvector for S associated with μ , we then have that $v \in E_{\mu} \subseteq M$, so $v \in M \cap M^{\perp} = \{0\}$, contradicting that $v \neq 0$ (since v is an eigenvector).

This means that S has to be 0. Thus we get T(y) = S(y) = 0 for all $y \in M^{\perp}$, as desired.

(e): If $\ker(T) = \{0\}$, then we get from (d) that $\mathcal{E} = \mathcal{E}'$ is an orthonormal basis for $\ker(T)^{\perp} = \{0\}^{\perp} = H$. If $\ker(T) \neq \{0\}$, then we have $\mathcal{E}_0 \subseteq \ker(T)$

and $\mathcal{E}' \subseteq \ker(T)^{\perp}$, so it is clear that \mathcal{E} is an orthonormal set. Moreover, we have that

$$H = \overline{\operatorname{Span}\left(\mathcal{E}\right)} \,.$$

Indeed, let $x \in H$. Then we may write

$$x = x_M + x_{M^\perp} \,,$$

where $x_M \in M = \overline{\text{Span}(\mathcal{E}')}$ and $x_{M^{\perp}} \in M^{\perp} = \text{ker}(T) = \overline{\text{Span}(\mathcal{E}_0)}$. So we may choose $\{x_n\}_{n\in\mathbb{N}} \subseteq \text{Span}(\mathcal{E}')$ and $\{y_n\}_{n\in\mathbb{N}} \subseteq \text{Span}(\mathcal{E}_0)$ such that $\lim_{n\to\infty} x_n = x_M$ and $\lim_{n\to\infty} y_n = x_{M^{\perp}}$. This gives that

$$\lim_{n \to \infty} (x_n + y_n) = x_M + x_{M^\perp} = x \,.$$

Hence, $x \in \text{Span}(\mathcal{E})$. This shows that \mathcal{E} is an orthonormal basis for H.

(f): The first assertion follows readily from the fact that $E_{\lambda} \perp E_{\lambda'}$ whenever $\lambda \neq \lambda'$, cf. Lemma 4.3.2. Next, we consider the case where L is countably infinite and $\{\lambda_k : k \in \mathbb{N}\}$ is an enumeration of L, leaving the easier case where L is finite to the reader.

For each $k \in \mathbb{N}$, set $n_k := \dim(E_{\lambda_k}) < \infty$, and let $\{v_{k,1}, \ldots, v_{k,n_k}\}$ be an enumeration of \mathcal{E}_{λ_k} . Then we have

$$\mathcal{E}' = \bigcup_{k \in \mathbb{N}} \mathcal{E}_{\lambda_k} = \{ v_{k,l} : k \in \mathbb{N}, 1 \le l \le n_k \}.$$

Consider $x \in H$. Since $T(x) \in T(H)$ and \mathcal{E}' is an orthonormal basis for $\overline{T(H)}$, we get from Corollary 3.3.11 that

$$T(x) = \lim_{m \to \infty} \sum_{k=1}^{m} \sum_{l=1}^{n_k} \langle T(x), v_{k,l} \rangle v_{k,l} = \lim_{m \to \infty} \sum_{k=1}^{m} \sum_{l=1}^{n_k} \langle x, T(v_{k,l}) \rangle v_{k,l}$$
$$= \lim_{m \to \infty} \sum_{k=1}^{m} \lambda_k \left(\sum_{l=1}^{n_k} \langle x, v_{k,l} \rangle v_{k,l} \right) = \sum_{k=1}^{\infty} \lambda_k P_{\lambda_k}(x).$$

Let now $\varepsilon > 0$. We have to show that there exists $N \in \mathbb{N}$ such that $||T - \sum_{k=1}^{n} \lambda_k P_{\lambda_k}|| \leq \varepsilon$ for all $n \geq N$.

Using (b), we can choose $N \in \mathbb{N}$ such that $|\lambda_k| < \varepsilon$ for all k > N. Then for all $n \ge N$ and all $x \in H$, using continuity of the norm in H and Pythagoras' identity, we get

$$\left\| \left(T - \sum_{k=1}^{n} \lambda_k P_{\lambda_k}\right)(x) \right\|^2 = \left\| \sum_{k=n+1}^{\infty} \lambda_k P_{\lambda_k}(x) \right\|^2 = \sum_{k=n+1}^{\infty} |\lambda_k|^2 \|P_{\lambda_k}(x)\|^2$$
$$\leq \varepsilon^2 \sum_{k=n+1}^{\infty} \|P_{\lambda_k}(x)\|^2 \leq \varepsilon^2 \|x\|^2$$

and the assertion follows.

Sometimes, the following consequence of the spectral theorem is useful.

Corollary 4.3.5. Let $T \in \mathcal{K}(H)$ be self-adjoint, $T \neq 0$. Then there exist an orthonormal family $\{u_j\}_{j\in J}$ in H and a family $\{\rho_j\}_{j\in J}$ in $\mathbb{R} \setminus \{0\}$, both indexed over the same countable set J, such that

$$T(x) = \sum_{j \in J} \rho_j \langle x, u_j \rangle u_j \quad \text{for all } x \in H.$$
(4.3.2)

Each ρ_i is an eigenvalue of T and u_i is an associated eigenvector.

Proof. We use the notation from Theorem 4.3.4). Since $\mathcal{E}' = \bigcup_{\lambda \in L} \mathcal{E}_{\lambda}$ is countable, we may let $\{u_j : j \in J\}$ be an enumeration of \mathcal{E}' , where the index set J is countable. Moreover, for each $j \in J$, we may then let $\rho_j \in \mathbb{R} \setminus \{0\}$ denote the eigenvalue of T corresponding to u_j . (The family $\{\rho_j\}_{j \in J}$ gives a list of all nonzero eigenvalues of T, repeated according to their multiplicites.) Since $\mathcal{E}' = \{u_j : j \in J\}$ is an orthonormal basis for $\overline{T(H)}$, as in the proof of (f), we get

$$T(x) = \sum_{j \in J} \langle T(x), u_j \rangle \, u_j = \sum_{j \in J} \langle x, T(u_j) \rangle \, u_j = \sum_{j \in J} \rho_j \, \langle x, u_j \rangle \, u_j$$

for all $x \in H$.

Let us say that an operator $T \in \mathcal{B}(H)$ is *diagonalizable* if there exists an orthonormal basis for H whose elements are eigenvectors for T. Then the spectral theorem (Theorem 4.3.4) says that T is diagonalizable whenever T is compact and self-adjoint. Let us say that an operator $T \in \mathcal{B}(H)$ is normal if T^* commutes with T, i.e., $T^*T = TT^*$. Clearly, self-adjoint operators and unitary operators are normal. A more precise version of Theorem 4.3.4 is as follows.

Theorem 4.3.6. Assume T is a compact operator on H. If $\mathbb{F} = \mathbb{R}$, then T is diagonalizable if and only if T is self-adjoint. On the other hand, if $\mathbb{F} = \mathbb{C}$, then T is diagonalizable if and only if T is normal.

This theorem can be deduced from Theorem 4.3.4. We leave the proof to the reader (cf. Exercises 4.13 and 4.14).

As another corollary of Theorem 4.3.4, we end this section by showing how an analogue of the singular value decomposition (SVD) for matrices may be obtained for compact operators.

Let $S \in \mathcal{K}(H)$, $S \neq 0$. Then $T := S^*S$ is self-adjoint and compact, and $T \neq 0$ (as $||T|| = ||S^*S|| = ||S||^2 \neq 0$). Hence, the spectral theorem (Theorem 4.3.4) gives that we can find a countable orthonormal basis $\{v_j\}_{j \in N}$ for

$$\overline{T(H)} = \ker(T)^{\perp} = \ker(S^*S)^{\perp} = \ker(S)^{\perp}$$

consisting of eigenvectors for T. For each $j \in N$, let μ_j denote the eigenvalue of T associated with v_j . Note that

$$\mu_j = \langle \mu_j \, v_j, v_j \rangle = \langle T(v_j), v_j \rangle = \langle S(v_j), S(v_j) \rangle = \|S(v_j)\|^2 \ge 0$$

for every $j \in N$. Since each μ_j is nonzero, we get that all μ_j 's are positive. For each $j \in N$, set

$$\sigma_j := \sqrt{\mu_j} > 0$$
 and $u_j := \frac{1}{\sigma_j} S(v_j)$.

The σ_j 's are called the singular values of S. For all $j, k \in N$ we have

$$\begin{aligned} \langle u_j, u_k \rangle &= \frac{1}{\sigma_j \sigma_k} \left\langle S(v_j), S(v_k) \right\rangle = \frac{1}{\sigma_j \sigma_k} \left\langle T(v_j), v_k \right\rangle \\ &= \frac{\mu_j}{\sigma_j \sigma_k} \left\langle v_j, v_k \right\rangle = \begin{cases} 1 & \text{if } j = k, \\ 0 & \text{otherwise,} \end{cases} \end{aligned}$$

so $\{u_j : j \in N\}$ is an orthonormal set in the range of S. Further, we have the following decomposition of S:

$$S(x) = \sum_{j \in N} \sigma_j \langle x, v_j \rangle u_j \quad \text{for all } x \in H.$$
(4.3.3)

Indeed, let $x \in H$ and set $M := \overline{T(H)}$, so $M^{\perp} = \ker(S)$. With $z := x - P_M(x) \in M^{\perp}$, we get that

$$x = P_M(x) + z = \sum_{j \in N} \langle x, v_j \rangle v_j + z,$$

SO

$$S(x) = \sum_{j \in N} \langle x, v_j \rangle S(v_j) + S(z) = \sum_{j \in N} \sigma_j \langle x, v_j \rangle u_j,$$

as asserted in (4.3.3).

It readily follows that $\{u_j : j \in N\}$ is an orthonormal basis for S(H). Finally we remark that the spectral theorem also gives that $\sigma_j = \sqrt{\mu_j} \to 0$ as $j \to \infty$ when N is countably infinite, and that

$$||S|| = ||S^*S||^{1/2} = ||T||^{1/2} = \max\{\mu_j : j \in N\}^{1/2} = \max\{\sigma_j : j \in N\}.$$

4.4 Application: The Fredholm Alternative

A useful application of linear algebra, and one of its original motivation, is the study of systems of linear equations, i.e., of equations of the type A x = b, where $A \in M_{m \times n}(\mathbb{F})$, $b \in \mathbb{F}^m$ and the (unknown) vector x belongs to \mathbb{F}^n . More generally, one may consider equations of the form

$$T(v) = w \tag{4.4.1}$$

where V, W are vector spaces (over \mathbb{F}), $T \in \mathcal{L}(V, W)$, $w \in W$ and the (unknown) vector x belongs to V. Whether such an equation is consistent, i.e., has some solution(s), relies on whether w lies in the range of T. If this is the case, and $v_0 \in V$ is any vector satisfying (4.4.1), i.e., such that $T(v_0) = w$, then it follows readily that the solution set of (4.4.1) is given by

$$v_0 + \ker(T) := \{ v_0 + u \mid u \in \ker(T) \}$$
 (4.4.2)

where $v_0 \in V$ is any vector satisfying (4.4.1), i.e., such that $T(v_0) = w$.

In the rest of this section, we consider the case where V = W = H is a Hilbert space ($\neq \{0\}$), and $T \in \mathcal{B}(H)$. We can then exploit the relationship between the fundamental subspaces of T and T^* , cf. Proposition 3.4.12.

For example, using that $\overline{T(H)} = \ker(T^*)^{\perp}$, we get that if T has closed range (i.e., T(H) is closed), then the equation (4.4.1) will be consistent if and only if w is orthogonal to $\ker(T^*)$.

In particular, if T has closed range and ker $(T^*) = \{0\}$ (i.e., T^* is one-toone), then T must be surjective, hence (4.4.1) is consistent for all $w \in H$. Similarly, if T^* has closed range and ker $(T) = \{0\}$, then it follows that T^* is surjective, so the equation $T^*(v') = w'$ is consistent for all $w' \in H$.

On the other hand, if it happens that T is surjective, then we get that $\ker(T^*) = \{0\}$, hence that the equation $T^*(v') = w'$ will have either no solution or a unique solution. Similarly, if T^* is surjective, then $\ker(T^*) = \{0\}$, and (4.4.1) will have either no solution or a unique solution.

A problem is that bounded operators in general do not have a closed range (see for example Exercises 4.17 and 4.18). Moreover, it may often be a difficult task to decide whether the range of some given $T \in \mathcal{B}(H)$ is closed or not. However, we note that if $T \in \mathcal{B}(H)$ has finite-rank, then it has closed range (as T(H) is finite-dimensional). In the case where H is finite-dimensional, much more can be said.

The following terminology will be useful.

Definition 4.4.1. An operator $F \in \mathcal{B}(H)$ is said to satisfy the *Fredholm* alternative if one of the following two (mutually exclusive) situations occurs:

- (a) $\ker(F) = \ker(F^*) = \{0\}$, and the equations F(v) = w, $F^*(v') = w'$ have both a unique solution for all $w, w' \in H$;
- (b) $1 \leq \dim(\ker(F)) = \dim(\ker(F^*)) < \infty$, the equation F(v) = w is consistent if and only if $w \in \ker(F^*)^{\perp}$, and the equation $F^*(v') = w'$ is consistent if and only if $w' \in \ker(F)^{\perp}$.

Example 4.4.2. Assume that *H* is finite-dimensional and $F \in \mathcal{B}(H)$, i.e., $F \in \mathcal{L}(H)$. Then *F* satisfies the Fredholm alternative.

The crux is that we have $\dim(\ker(F^*)) = \dim(\ker(F))$. To show this, we use the formula $\dim(M) + \dim(M^{\perp}) = \dim(H)$,

which is easily verified for any subspace M of H, and the dimension formula for F. We get that

 $\dim(\ker(F^*)) = \dim(F(H)^{\perp}) = \dim(H) - \dim(F(H)) = \dim(\ker(F)).$

Combining this fact with our previous observations, it is straightforward to deduce that either (a) or (b) in Definition 4.4.1 holds. \Box

When T_K is an integral operator on $L^2([a, b])$, and $\mu \in \mathbb{C}$, an equation of the form $(T_K - \mu I)(f) = g$, i.e., $T_K(f) - \mu f = g$, is often called a Fredholm integral equation of the second kind. Such equations, and Fredholm integral equations of the first kind (i.e., equations of the form $T_K(f) = g$), were studied by I. Fredholm at the beginning of the 20th century. They arise in some practical problems in signal theory and in physics. He showed that if $\mu \neq 0$, then the operator $T_K - \mu I$ satisfies the Fredholm alternative. Now, as we have seen in Example 4.2.11, T_K is compact operator. It can in fact be shown that if T is a compact operator on H and $\mu \in \mathbb{F} \setminus \{0\}$, then any operator of the form $T - \mu I$ satisfies the Fredholm alternative.

To give an idea of the proof of this result, consider $T \in \mathcal{K}(H)$ and $\mu \in \mathbb{F} \setminus \{0\}$. Then it can be shown that the following facts hold:

- (i) $T \mu I$ has closed range;
- (ii) $\dim(\ker(T-\mu I)) = \dim(\ker((T-\mu I)^*)) < \infty.$

Since T^* is compact, we also get that $T^* - \overline{\mu}I = (T - \mu I)^*$ has closed range. Using these properties, and the general principles outlined before, one readily arrives at the conclusion that $T - \mu I$ satisfies the Fredholm alternative, as asserted above. We don't have time in this course to prove that (i) and (ii) hold. Instead, we will illustrate how the spectral theorem for compact self-adjoint operators can be applied to give a direct proof of the following result.

Theorem 4.4.3. Assume $T \in \mathcal{K}(H)$ is self-adjoint and $\mu \in \mathbb{F} \setminus \{0\}$. Then $F = T - \mu I$ satisfies the Fredholm alternative.

Proof. Assume first that μ is not an eigenvalue of T, i.e., $\ker(T - \mu I) = \{0\}$. Then the spectral theorem implies that the equation $(T - \mu I)(x) = y$ has a unique solution for all $y \in H$. (You are asked to check this in Exercise 4.9.) Thus, $F = T - \mu I$ is surjective, and this implies that $\ker(F^*) = \ker(T - \overline{\mu}I) = \{0\}$, i.e., $\overline{\mu}$ is not an eigenvalue of T. Arguing as above, we get that the equation $(T - \overline{\mu}I)(x') = y'$, i.e., $(T - \mu I)^*(x') = y'$ has a unique solution for all $y' \in H$. This shows that (a) in Definition 4.4.1 holds in this case.

Next, assume that μ is an eigenvalue of T, i.e., $\ker(T - \mu I) \neq \{0\}$. Then $\mu \in \mathbb{R}$, so $F^* = F$. Moreover, as $\mu \neq 0$, we have that $T \neq 0$, and the spectral theorem tells us that $1 \leq \dim(\ker(F)) = \dim(\ker(T - \mu I)) < \infty$. Hence, to show that (b) in Definition 4.4.1 holds, it remains only to prove that the equation F(x) = y is consistent if and only if $y \in \ker(F)^{\perp}$. This means that we have to prove that the equation

$$T(x) - \mu x = y \tag{4.4.3}$$

is consistent if and only if $\langle y, z \rangle = 0$ for all $z \in E_{\mu} := \ker(T - \mu I)$.

To prove this, let $\mathcal{E}' = \{u_j\}_{j \in J}$ be an enumeration of the orthonormal basis for $\overline{T(H)}$ obtained in the spectral theorem for T, and let $\rho_j \in \mathbb{R} \setminus \{0\}$ denote the eigenvalue of T corresponding to each u_j (as in the proof of Corollary 4.3.5).

Since H is the direct sum of $\ker(T)$ and $\ker(T)^{\perp} = \overline{T(H)}$, we may write $y \in H$ as

$$y = y_0 + \sum_{j \in J} \langle y, u_j \rangle \, u_j,$$

where y_0 denote the orthogonal projection of y onto ker(T). Likewise, we may assume that the (unknown) vector x in equation (4.4.3) is written as

$$x = x_0 + \sum_{j \in J} c_j \, u_j,$$

where $x_0 \in \ker(T)$ and $\{c_j\}_{j \in J} \in \ell^2(J)$ are to be determined, if possible. Plugging this into equation (4.4.3), we get the equivalent equation

$$-\mu x_0 + \sum_{j \in J} \left(\rho_j - \mu \right) c_j \, u_j = y_0 + \sum_{j \in J} \left\langle y, u_j \right\rangle u_j.$$

Clearly, we can set $x_0 := (-1/\mu) y_0$, and equation (4.4.3) is then consistent if and only if the sequence $\{c_j\}_{j \in J} \in \ell^2(J)$ can be chosen so that

$$(\rho_j - \mu) c_j = \langle y, u_j \rangle \text{ for all } j \in J.$$

$$(4.4.4)$$

Now, as μ is a nonzero eigenvalue of T, we have that $\mu = \rho_k$ for some $k \in J$. Let u_{j_1}, \ldots, u_{j_n} denote the vectors in \mathcal{E}' giving an orthonormal basis for $E_{\mu} = E_{\rho_k}$. If $j \notin \{j_1, \ldots, j_n\}$, we have $\rho_j \neq \mu$, so

$$c_j := \frac{1}{\rho_j - \mu} \left\langle y, u_j \right\rangle$$

will satisfy (4.4.4) for every such j.

On the other hand, if $j \in \{j_1, \ldots, j_n\}$, we have $\rho_j - \mu = 0$. Hence, (4.4.4) will be satisfied for $j = j_1, \ldots, j_n$ if and only if we have $\langle y, u_j \rangle = 0$ for $j = j_1, \ldots, j_n$, i.e., if and only if $\langle y, z \rangle = 0$ for all $z \in E_{\mu}$. Moreover, when this condition holds, we can choose c_{j_1}, \ldots, c_{j_n} freely and, regardless of this choice, the constructed sequence $\{c_j\}_{j \in J}$ is easily seen to belong to $\ell^2(J)$ (cf. Exercise 4.20), meaning that the associated vector x gives a solution to (4.4.3). Thus, we have proved the desired equivalence.

4.5 Exercises

Exercise 4.1. Let X, Y, Z denote normed spaces over \mathbb{F} . Consider $\lambda \in \mathbb{F}$, $T, T' \in \mathcal{B}(X, Y)$ and $S \in \mathcal{B}(Y, Z)$, so $ST \in \mathcal{B}(X, Z)$.

- a) Show that $\lambda T + T' \in \mathcal{K}(X, Y)$ if $T, T' \in \mathcal{K}(X, Y)$.
- b) Show that $ST \in \mathcal{K}(X, Z)$ if $T \in \mathcal{K}(X, Y)$.
- c) Show that $ST \in \mathcal{K}(X, Z)$ if $S \in \mathcal{K}(Y, Z)$.
- d) Set $\mathcal{K}(X) = \mathcal{K}(X, X)$. Deduce that
- $ST \in \mathcal{K}(X)$ if $S \in \mathcal{B}(X)$ and $T \in \mathcal{K}(X)$, or if $S \in \mathcal{K}(X)$ and $T \in \mathcal{B}(X)$.

Exercise 4.2. Let $X = \ell^p(\mathbb{N}), \lambda \in \ell^\infty(\mathbb{N})$, and $M_\lambda \in \mathcal{B}(X)$ be the associated multiplication operator, cf. Example 4.1.8.

Show that $\lambda \in c_0(\mathbb{N})$ if M_{λ} is compact.

(It therefore follows that M_{λ} is compact if and only if $\lambda \in c_0(\mathbb{N})$.)

Exercise 4.3. Let X be a normed space, H be a Hilbert space, and let $T \in \mathcal{K}(X, H)$. Show that $\overline{T(X)}$ is separable.

Exercise 4.4. Let H be an infinite-dimensional Hilbert space and let $T \in \mathcal{K}(H)$. Show that $\langle T(u_n), u_n \rangle \to 0$ as $n \to \infty$ whenever $\{u_n\}_{n \in \mathbb{N}}$ is an orthonormal sequence in H.

Exercise 4.5. Let *H* be a Hilbert space and assume $P \in \mathcal{B}(H)$ satisfies $P^2 = P$.

Show that P has finite-rank if (and only if) P is compact.

Exercise 4.6. Let *H* be a separable Hilbert space, $H \neq \{0\}$.

a) Show that $\mathcal{F}(H) \subseteq \mathcal{HS}(H)$, and that $\mathcal{F}(H)$ is dense in $\mathcal{HS}(H)$ w.r.t. $\|\cdot\|_2$.

b) Assume that $T \in \mathcal{HS}(H)$ and $S \in \mathcal{B}(H)$. Show that both ST and TS belong to $\mathcal{HS}(H)$, and that we have

$$||ST||_2 \le ||S|| ||T||_2, ||TS|| \le ||T||_2 ||S||.$$

c) Let $\mathcal{B} = \{u_j\}_{j \in J}$ be an orthonormal basis for H, where $J = \{1, \ldots, n\}$ if $\dim(H) = n < \infty$, while $J = \mathbb{N}$ otherwise.

For $T, T' \in \mathcal{HS}(H)$, set

$$\langle T, T' \rangle_2 := \sum_{j \in J} \langle T(u_j), T'(u_j) \rangle$$

Show that this gives a well-defined inner product on $\mathcal{HS}(H)$, and check that the associated norm is the Hilbert-Schmidt norm $\|\cdot\|_2$.

d) Show that $\mathcal{HS}(H)$ is complete w.r.t. $\|\cdot\|_2$, so that $\mathcal{HS}(H)$ is a Hilbert space w.r.t. the inner product from c).

Exercise 4.7. Let $H = L^2(\mathbb{R}, \mathcal{A}, \mu)$ where \mathcal{A} denote all Lebesgue measurable subsets of \mathbb{R} and μ is the Lebesgue measure. For which $f \in \mathcal{L}^{\infty}$ is the multiplication operator $M_f \in \mathcal{B}(H)$ compact?

Exercise 4.8. Let H be a Hilbert space, $T \in \mathcal{K}(H)$ and $\lambda \in \mathbb{F}, \lambda \neq 0$. Assume that there exists a sequence $\{x_n\}_{n\in\mathbb{N}}$ of unit vectors in H such that $||T(x_n) - \lambda x_n|| \to 0$ as $n \to \infty$. Show that λ is an eigenvalue of T.

Exercise 4.9. Let H be a Hilbert space, and let $T \in \mathcal{K}(H)$ be self-adjoint. Assume $\mu \in \mathbb{F}, \mu \neq 0$ is *not* an eigenvalue of T, i.e. $T - \mu I_H$ is injective.

Let $y \in H$, let $\mathcal{E}' = \{u_j\}_{j \in J}$ be an enumeration of the orthonormal basis for $M = \overline{T(H)}$ obtained in the spectral theorem for T, and let $\rho_j \neq 0$ denote the eigenvalue of T corresponding to u_j (as in the proof of Corollary 4.3.5). a) Show that the series

$$\sum_{j \in J} \frac{\langle y, u_j \rangle}{\rho_j - \mu} u_j$$

converges to some $h \in H$.

- b) Set $z := y P_M(y)$ and $x := h \frac{1}{\mu} z$. Show that $(T \mu I_H)(x) = y$.
- c) Deduce that $T \mu I_H$ is surjective (hence that it is bijective).

Exercise 4.10. Consider $H = L^2([-\pi, \pi])$ (with respect to the normalized Lebesgue measure). Let $g \in C([-\pi, \pi])$ be periodic, i.e., satisfies that $g(-\pi) = g(\pi)$, and extend g to a periodic function \tilde{g} on \mathbb{R} with period 2π . Define $G : [-\pi, \pi] \times [-\pi, \pi] \to \mathbb{C}$ by

$$G(s,t) = \tilde{g}(s-t) \,.$$

a) Check that G is continuous, so that the associated integral operator T_G belongs to $\mathcal{HS}(H)$ (hence is compact).

c) Decide when T_G is self-adjoint.

b) Let $k \in \mathbb{Z}$ and recall that $e_k(t) = e^{ikt}$ for all $t \in [-\pi, \pi]$. Check that e_k is an eigenvector for the operator T_G . Deduce that T_G is diagonalizable (with respect to $\{e_k\}_{k\in\mathbb{Z}}$).

c) Show that $||T_G||_2 = ||g||_2 = \left(\frac{1}{2\pi} \int_{-\pi}^{\pi} |g(t)|^2 dt\right)^{1/2}$.

Exercise 4.11. Consider $H = L^2([0, 1])$ (with respect to Lebesgue measure) and the integral operator $T_K \in \mathcal{B}(H)$ associated to the kernel given by

$$K(s,t) = \min(s,t)$$

for all (s, t) in $[0, 1] \times [0, 1]$, cf. Example 4.2.11.

a) Explain why T_K is self-adjoint and compact. Then check that the set $\mathcal{U} := \{[u_n] : n \in \mathbb{N}\},$ where

$$u_n(t) := \sqrt{2} \sin\left(\left(n - \frac{1}{2}\right)\pi t\right) \quad \text{for all } t \in [0, 1], \ n \in \mathbb{N},$$

is an orthonormal set of eigenvectors for T_K .

b) It can be shown that \mathcal{U} is an orthonormal basis for H. Is it possible to deduce this from a) and the spectral theorem for T_K ?

Exercise 4.12. Let $S, T \in \mathcal{B}(H)$.

a) Assume there exists an orthonormal basis for H whose elements are eigenvectors for both S and T. Check that S commutes with T.

b) Assume S and T are compact and self-adjoint, and that S commutes with T. Show that there exists an orthonormal basis for H whose elements are eigenvectors for both S and T.

Hint: Start by considering an eigenvalue λ of T and study how S acts on the corresponding eigenspace E_{λ}^{T} .

Exercise 4.13. Assume H is a Hilbert space over \mathbb{R} , and let $T \in \mathcal{B}(H)$.

a) Assume that T is diagonalizable (as defined in Remark 4.3.6). Check that T is self-adjoint.

b) Let T be compact. Deduce that T is diagonalizable if and only if T is self-adjoint.

Exercise 4.14. Assume H is a Hilbert space over \mathbb{C} , and let $T \in \mathcal{B}(H)$.

a) Assume that T is diagonalizable (as defined in Remark 4.3.6). Check that T is normal.

b) Show that T is normal if and only if $\operatorname{Re}(T)$ and $\operatorname{Im}(T)$ commutes with each other.

c) Let T be compact. Show that T is diagonalizable if and only if T is normal.

Hint: The implication (\Rightarrow) follows from a). For (\Leftarrow) , use b) and Exercise 4.12 b).

Exercise 4.15. Let H be a separable Hilbert space with a countably infinite orthonormal basis $\mathcal{B} = \{v_j\}_{j \in \mathbb{N}}$. Let $\{\mu_j\}_{j \in \mathbb{N}}$ be a bounded sequence in \mathbb{F} and let $D \in \mathcal{B}(H)$ denote the associated diagonal operator (w.r.t. \mathcal{B}).

a) Show that D is compact if and only if $\lim_{j\to\infty} \mu_j = 0$.

(*Note*: If you have looked at Example 4.1.8 and solved Exercise 4.2, this should not be difficult).

b) Show that D is Hilbert-Schmidt if and only if $\{\mu_j\}_{j\in\mathbb{N}} \in \ell^2(\mathbb{N})$, in which case we have

$$||D||_2 = \left(\sum_{j=1}^{\infty} |\mu_j|^2\right)^{1/2}.$$

Exercise 4.16. Let H be a separable Hilbert space of infinite dimension and let $T \in \mathcal{K}(H)$ be selfadjoint, $T \neq 0$. Assume that you have found an orthonormal basis $\mathcal{B} = \{v_j\}_{j \in \mathbb{N}}$ for H consisting of eigenvectors for T, and let $\mu_j \in \mathbb{R}$ denote the eigenvalue of T corresponding to each v_j .

a) Show that the sequence $\{\mu_j\}_{j\in\mathbb{N}}$ is bounded, hence that T is the diagonal operator (w.r.t. \mathcal{B}) associated with this sequence. Deduce from the previous exercise that $\lim_{j\to\infty} \mu_j = 0$.

b) As in the spectral theorem, set

 $L := \{ \lambda \in \mathbb{R} \mid \lambda \text{ is a nonzero eigenvalue of } T \}.$

Set also

 $\widetilde{L} := \{ \lambda \in \mathbb{R} \mid \lambda \text{ is an eigenvalue of } T \},\$

so $L = L \setminus \{0\}$. Show the following assertions:

- (i) $\widetilde{L} = \{\mu_j \mid j \in \mathbb{N}\}$ and $L = \{\mu_j \mid j \in \mathbb{N}, \mu_j \neq 0\}.$
- (ii) If $\lambda \in L$ and $N_{\lambda} := \{j \in \mathbb{N} \mid \mu_j = \lambda\}$, then N_{λ} is a finite subset of \mathbb{N} and $\{v_j \mid j \in N_{\lambda}\}$ is an o.n.b. for E_{λ} .
- (iii) If $\mu_j \neq 0$ for all $j \in \mathbb{N}$, then ker $(T) = \{0\}$.
- (iv) If $N_0 := \{j \in \mathbb{N} \mid \mu_j = 0\}$ is nonempty, then $\{v_j \mid j \in N_0\}$ is an o.n.b. for ker(T).

Exercise 4.17. Let $H = L^2([0, 1])$ (with usual Lebesgue measure) and let $T = M_f$ be the self-adjoint operator in $\mathcal{B}(H)$ given by multiplication with the function f(t) = t on [0, 1], cf. Example 3.5.4.

Show that T(H) is not closed, i.e., that T does not have closed range. Show also that T is not compact.

Exercise 4.18. Let $H = \ell^2(\mathbb{N})$, let $\lambda \in \ell^{\infty}(\mathbb{N})$ be given by $\lambda(n) = \frac{1}{n}$ for all $n \in \mathbb{N}$, and let $T = M_{\lambda} \in \mathcal{B}(H)$ denote the associated multiplication operator. Note that T is compact, as follows from Example 4.1.8.

Show that T(H) = H and $T(H) \neq H$, so T does not have closed range.

Exercise 4.19. Let H be a Hilbert space and $T \in \mathcal{B}(H)$. Let us say that T is *bounded from below* if there exists some $\alpha > 0$ such that $\alpha ||x|| \le ||T(x)||$ for all $x \in H$. For example, T is bounded from below when T is an isometry.

Show that if T is bounded from below, then T has closed range.

Exercise 4.20. Finish the proof of Theorem 4.4.3 by checking that the sequence $\{c_j\}_{j\in J}$ constructed in the final paragraph (under the assumption that y is orthogonal to E_{μ}) belongs to $\ell^2(J)$.

CHAPTER 5

An introduction to Sturm-Liouville theory

The purpose of this chapter is to illustrate how the spectral theorem for a compact self-adjoint operator on a Hilbert space may be used to study some classical Sturm-Liouville problems. For simplicity we will only discuss the so-called regular case.

5.1 Regular Sturm-Liouville systems

When $[a, b] \subset \mathbb{R}$ and $n \in \mathbb{N}$, we let $C^n([a, b])$ denote the space of n times continuously differentiable complex functions on [a, b]. A regular Sturm-Liouville system on [a, b] is a second order linear differential equation of the form

$$-(py')' + q y = \lambda \rho y, \qquad (5.1.1)$$

where

- $p \in C^1([a, b])$ is real-valued and $p(x) \neq 0$ for all $x \in [a, b]$,
- $q, \rho \in C([a, b])$ are real-valued and $\rho(x) \neq 0$ for all $x \in [a, b]$,
- $\lambda \in \mathbb{C}$,

and the unknown function y = y(x), which necessarily has to lie in $C^2([a, b])$, is required to satisfy boundary conditions of the type

$$\alpha_1 y(a) + \alpha_2 y'(a) = 0, \quad \beta_1 y(b) + \beta_2 y'(b) = 0, \quad (5.1.2)$$

for some $(\alpha_1, \alpha_2), (\beta_1, \beta_2) \in \mathbb{R}^2 \setminus \{(0, 0)\}.$

5. An introduction to Sturm-Liouville theory

Ideally, the Sturm-Liouville problem is to determine the values of λ for which there exist non-trivial solutions of equation (5.1.1) satisfying the conditions (5.1.2), and to describe these solutions. These values of λ are called the *eigenvalues* of the system, and the corresponding solutions y are called *eigenfunctions* of the system. A concrete answer to this problem is not possible in general, but as we will see, one may still obtain some valuable theoretical information about it.

Since we only intend to give a small taste of Sturm-Liouville theory, we will assume that $p(x) = \rho(x) = 1$ for all $x \in [a, b]$, in which case equation (5.1.1) simplifies to

$$-y'' + qy = \lambda y. \tag{5.1.3}$$

A suitably scaled version of this equation appears for example as the one dimensional time-independent Schrödinger equation in quantum mechanics (where it is usually considered on the whole real line).

Set $Y = \{y \in C^2([a, b]) : y \text{ satisfies the boundary conditions } (5.1.2)\}$. We will consider Y and C([a, b]) as inner product spaces w.r.t. to the inner product given by

$$\langle f,g\rangle = \int_{a}^{b} f(x)\overline{g(x)} \, dx.$$

Letting $D: Y \to C[a, b]$ be the linear operator defined by

$$D(y) = -y'' + qy,$$

it is clear that our Sturm-Liouville system may be written as

$$D(y) = \lambda y$$
 where $y \in Y$. (5.1.4)

Although the associated Sturm-Liouville problem looks like a familiar eigenvalue/eigenvector problem, it is not obvious how to proceed. The fact that Y and C([a, b]) are not Hilbert spaces (they are not complete) can easily be fixed because both can be considered as dense subspaces of $L^2([a, b])$. (We leave it as an exercise to show this for Y). However, the trouble is that D is not a bounded operator (check this!), so it does not extend to a bounded operator on $L^2([a, b])$. We will have to work quite a bit to recast the problem into one involving a compact selfadjoint operator.

• We will first study the second order differential equation

$$-y'' + qy = \lambda y \quad \text{with } y \in C^2([a, b]) \tag{5.1.5}$$

and show that its solution space

$$S_{\lambda} := \{ y \in C^2([a,b]) : -y'' + qy = \lambda y \}$$

is 2-dimensional for every $\lambda \in \mathbb{C}^{1}$ Note that trying to find out when there exists some $y \in S_{\lambda} \setminus \{0\}$ which also belongs to Y, which would solve our problem, is not possible because a concrete description of S_{λ} is not available in general.

- Next, we will establish some spectral properties of the operator D.
- Thirdly, we will assume that D is 1-1. We will then show that D is onto C([a, b]), and that there exists a compact self-adjoint operator $T_G: L^2([a, b]) \to L^2([a, b])$ such that its restriction to C([a, b]) is the inverse of D. Applying the spectral theorem to T_G will lead us to a theoretical answer to our Sturm-Liouville problem in this case.
- Finally, we will explain how to handle the general case where D is not assumed to be 1-1.

5.2 A second order differential equation

We recall that $\lambda \in \mathbb{C}$. In this section it is not important that the function $q \in C([a, b])$ is assumed to be real-valued.

Theorem 5.2.1. Let $c \in [a, b]$ and $z_1, z_2 \in \mathbb{C}$. Then there exists a unique function $y \in C^2([a, b])$ satisfying that

$$-y'' + qy = \lambda y$$
 and $y(c) = z_1, y'(c) = z_2.$ (5.2.1)

Proof. Suppose first that $y \in C^2([a, b])$ satisfies (5.2.1), that is,

$$y'' = (\lambda - q)y, \ y(c) = z_1 \text{ and } y'(c) = z_2.$$
 (5.2.2)

For every $u \in [a, b]$, we get

$$y'(u) - z_2 = y'(u) - y'(c) = \int_c^u y''(t) dt = \int_c^u (\lambda - q(t)) y(t) dt$$

¹This fact holds for the solution space of any homogeneous second order linear ordinary differential equation, as some students may have seen in a previous course. We will give a self-contained proof in our case.

This gives that

$$y(x) - z_1 - z_2(x - c) = y(x) - y(c) - z_2(x - c)$$

= $\int_c^x (y'(u) - z_2) du$
= $\int_c^x \int_c^u (\lambda - q(t)) y(t) dt du$
= $\int_c^x \int_t^x (\lambda - q(t)) y(t) du dt$
= $\int_c^x (x - t)(\lambda - q(t)) y(t) dt$

for all $x \in [a, b]$, hence that y satisfies that the integral equation

$$y(x) = z_1 + z_2(x-c) + \int_c^x (x-t)(\lambda - q(t)) y(t) dt \quad \text{for all } x \in [a, b].$$
(5.2.3)

Conversely, if $y \in C([a, b])$ satisfies (5.2.3), then it is an easy exercise to check that y belongs to $C^2([a, b])$ and satisfies (5.2.2).

Now, let $T: C([a, b]) \to C([a, b])$ be the integral operator defined for each f in C([a, b]) by

$$[T(f)](x) = z_1 + z_2(x - c) + \int_c^x (x - t)(\lambda - q(t)) f(t) dt$$

for all $x \in [a, b]$. We consider here C([a, b]) as a complete metric space w.r.t. the metric $d(f, g) := ||f - g||_{\infty} = \sup\{|f(x) - g(x)| : x \in [a, b]\}.$

Set $K := \sup \{ |(x-t)(\lambda - q(t))| : x, t \in [a, b] \} < \infty$. Let $f, g \in C([a, b])$ and $x \in [a, b]$. By induction on $n \in \mathbb{N}$, one easily shows that

$$\left| [T^{n}(f) - T^{n}(g)](x) \right| \leq \frac{1}{n!} K^{n} |x - c|^{n} ||f - g||_{\infty}.$$

This implies that

$$||T^{n}(f) - T^{n}(g)||_{\infty} \le \frac{K^{n} (b-a)^{n}}{n!} ||f - g||_{\infty}$$

It clearly follows that T^n is a contraction when n is so large that $\frac{K^n (b-a)^n}{n!} < 1$. Hence, Banach's fixed point theorem ² gives that T has a unique fixed point, say y, in C([a, b]). This means that y is the unique function in C([a, b]) such that y = T(y), i.e., such that y satisfies (5.2.3). Taking into account what we proved in the first part of the proof, we are done.

 $^{^{2}}$ cf. Lindstrøm's book *Spaces*, Exercise 3.4.7.
Corollary 5.2.2. Let $\lambda \in \mathbb{C}$, $c \in [a, b]$, and recall that

$$S_{\lambda} = \{ y \in C^2([a, b]) : -y'' + qy = \lambda y \}.$$

Then the map $T_{\lambda,c}: S_{\lambda} \to \mathbb{C}^2$ defined by $T_{\lambda,c}(y) = (y(c), y'(c))$ for every $y \in S_{\lambda}$ is an isomorphism. Hence, dim $S_{\lambda} = 2$.

Proof. Theorem 5.2.1 shows that the map $T_{\lambda,c}$ is 1-1 and onto. It is obvious that it is linear.

Remark 5.2.3. It should be noted that Theorem 5.2.1 is essentially an existence result (although our method of proof gives a way to approximate the unique solution of (5.2.1) by picking some $y_0 \in C([a, b])$ and computing $T^n(y_0)$ for large enough n). Explicit formulas for a basis of S_{λ} are only known when q is a constant function. To illustrate Corollary 5.2.2, we recall these. Assume $q(x) = \omega$ for all $x \in [a, b]$ for some $\omega \in \mathbb{C}$. Then $-y'' + qy = \lambda y$ can be rewritten as the homogeneous equation $y'' + (\lambda - \omega)y = 0$, which we know can be solved by considering the characteristic equation $z^2 + (\lambda - \omega) = 0$:

If $\lambda \neq \omega$, then, letting $(\omega - \lambda)^{1/2}$ denote a square root of $\omega - \lambda$ in \mathbb{C} , we get that S_{λ} consists of the functions of the form

$$y(x) = C_1 e^{(\omega - \lambda)^{1/2}x} + C_2 e^{-(\omega - \lambda)^{1/2}x}, \quad x \in [a, b],$$

where $C_1, C_2 \in \mathbb{C}$. Thus $\{e^{(\omega-\lambda)^{1/2}x}, e^{-(\omega-\lambda)^{1/2}x}\}$ is a basis for S_{λ} in this case.

If $\lambda = \omega$, the equation is y'' = 0 and $\{1, x\}$ is obviously a basis for S_{ω} . \Box

Remark 5.2.4. Suppose that q is real-valued, $\lambda \in \mathbb{R}$ and $y \in S_{\lambda}$. Then it is not difficult to verify that $\overline{y} \in S_{\lambda}$, so that $\operatorname{Re} y$ and $\operatorname{Im} y$ also lie in S_{λ} . Moreover, if it happens that y(c) and y'(c) both are real numbers for some $c \in [a, b]$, then the function y has to be real-valued: indeed, we then have $(\operatorname{Im} y)(c) = 0 = (\operatorname{Im} y)'(c)$, so Theorem 5.2.1 implies that $\operatorname{Im} y$ is the zero function on [a, b].

Remark 5.2.5. Consider $y_1, y_2 \in C^2([a, b])$. Define $W_{y_1, y_2} \in C^1([a, b])$ by

$$W_{y_1,y_2}(x) = \begin{vmatrix} y_1(x) & y_2(x) \\ y'_1(x) & y'_2(x) \end{vmatrix} \quad \text{for each } x \in [a,b].$$

 $W_{y_1,y_2}(x)$ is called the Wronsky determinant of (y_1, y_2) at x.

Assume that $y_1, y_2 \in S_{\lambda}$ and let $c \in [a, b]$. Corollary 5.2.2 implies that the set $\{y_1, y_2\}$ is a basis for S_{λ} if and only if the vectors $(y_1(c), y'_1(c)), (y_2(c), y'_2(c))$ are linearly independent in \mathbb{C}^2 , i.e., $W_{y_1,y_2}(c) \neq 0$. Note that this gives that if $\{y_1, y_2\}$ is a basis for S_{λ} , then $W_{y_1,y_2}(x) \neq 0$ for all $x \in [a, b]$.

The Wronsky determinant appears in the following lemma (sometimes called *Lagrange's lemma*), which will be useful to us later:

Lemma 5.2.6. Define $\widetilde{D}: C^2([a,b]) \to C([a,b])$ by

 $\widetilde{D}(y) = -y'' + qy \quad for \ every \ y \in C^2([a, b]),$

and let $f, g \in C^2([a, b])$. Then the following identity holds:

$$D(f)g - D(g)f = (fg' - gf')' = (W_{f,g})'.$$

Proof. We have

$$(W_{f,g})' = (fg' - gf')' = fg'' + f'g' - gf'' - g'f' = fg'' - gf'' = -f''g + qfg - qgf + fg'' = \widetilde{D}(f)g - \widetilde{D}(g)f.$$

5.3 Some spectral properties

Let $(\alpha_1, \alpha_2), (\beta_1, \beta_2) \in \mathbb{R}^2 \setminus \{(0, 0)\}$ and $q \in C([a, b])$ be real-valued. We recall that

 $Y = \Big\{ y \in C^2([a,b]) \mid y \text{ satisfies the boundary conditions } (5.3.1) \Big\},$ where

$$\alpha_1 y(a) + \alpha_2 y'(a) = 0, \quad \beta_1 y(b) + \beta_2 y'(b) = 0, \quad (5.3.1)$$

and $D: Y \to C[a, b]$ is the operator defined by D(y) = -y'' + qy for $y \in Y$.

Let $\lambda \in \mathbb{C}$ and set $E_{\lambda} := \{y \in Y : D(y) = \lambda y\}$. We say that λ is an eigenvalue of D if the subspace E_{λ} is non-trivial, in which case E_{λ} is called the eigenspace of D associated to λ . We note that $E_{\lambda} \subset S_{\lambda}$, so Corollary 5.2.2 implies that dim $E_{\lambda} \leq 2$.

Proposition 5.3.1. Let $f, g \in Y$. Then we have

- i) D(f)g D(g)f = (fg' gf')',
- *ii*) $\langle D(f), g \rangle = \langle f, D(g) \rangle$.

Proof. i) Since $D = \widetilde{D}_{|Y}$, this identity follows from Lagrange's lemma (Lemma 5.2.6). ii) It is easy to check that $\overline{g} \in Y$ and $\overline{D(g)} = D(\overline{g})$. Thus, using i), we get

$$\begin{split} \langle D(f),g\rangle - \langle f,D(g)\rangle &= \int_a^b [D(f)\overline{g} - fD(\overline{g})](t) \, dt \\ &= \int_a^b (f\overline{g}' - \overline{g}f')'(t) \, dt = \left[(f\overline{g}' - \overline{g}f')(t) \right]_a^b \\ &= f(b)\overline{g'(b)} - \overline{g(b)}f'(b) - f(a)\overline{g'(a)} + \overline{g(a)}f'(a). \end{split}$$

Now, since f and \overline{g} both satisfy 5.3.1, we have

$$\left[\begin{array}{cc} f(b) & f'(b) \\ \overline{g}(b) & \overline{g}'(b) \end{array}\right] \left[\begin{array}{c} \beta_1 \\ \beta_2 \end{array}\right] = \left[\begin{array}{c} 0 \\ 0 \end{array}\right].$$

Since $(\beta_1, \beta_2) \neq (0, 0)$, this implies that $f(b)\overline{g'(b)} - \overline{g(b)}f'(b) = 0$.

Arguing in a similar way, one can also show that $f(a)\overline{g'(a)} - \overline{g(a)}f'(a) = 0$. Inserting these two equalities in our computation above, we get that

$$\langle D(f), g \rangle - \langle f, D(g) \rangle = 0$$

as desired.

Part ii) of Proposition 5.3.1 shows that the operator D enjoys a property similar to self-adjointness. Proceeding exactly as we did for bounded self-adjoint operators on Hilbert spaces, one deduces that the following result holds.

Corollary 5.3.2. All the possible eigenvalues of D are real, and the associated eigenspaces are orthogonal to each other.

Note that we don't know yet whether D has any eigenvalues. Anyhow, we can say more about its eigenspaces (if any).

Proposition 5.3.3. All possible eigenspaces of D are one-dimensional.

Proof. Let $\lambda \in \mathbb{C}$. Recall that $D : C^2([a,b]) \to C([a,b])$ is defined by $\widetilde{D}(y) = -y'' + qy$ for $y \in C^2([a,b])$, so $D = \widetilde{D}_{|Y}$. We first consider the space

$$L_{\lambda} := \{ y \in C^2([a, b]) : D(y) = \lambda y \text{ and } y \text{ satisfies } (5.3.2) \},\$$

where

$$\alpha_1 \, y(a) + \alpha_2 \, y'(a) = 0. \tag{5.3.2}$$

Note that the condition (5.3.2) says that the vector (y(a), y'(a)) belongs to $M := \text{Span}\{(-\alpha_2, \alpha_1)\}$. Now, Corollary 5.2.2 (with c = a) gives that $L_{\lambda} = T_{\lambda,a}^{-1}(M)$. Since dim M = 1 and $T_{\lambda,a}$ is an isomorphism, we get that L_{λ} is a one-dimensional subspace of S_{λ} .

Similarly, one shows that

$$R_{\lambda} := \{ y \in C^2([a, b]) : \overline{D}(y) = \lambda y \text{ and } y \text{ satisfies } (5.3.3) \},\$$

where

$$\beta_1 y(b) + \beta_2 y'(b) = 0, \qquad (5.3.3)$$

is also a one-dimensional subspace of S_{λ} .

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Clearly, we have $E_{\lambda} = L_{\lambda} \cap R_{\lambda}$. So there are only two possibilities: either $E_{\lambda} = \{0\}$ or dim $E_{\lambda} = 1$. Hence, if λ is an eigenvalue of D, we must have dim $E_{\lambda} = 1$.

Remark 5.3.4. We use the notation introduced in the proof above. Assume that $\lambda \in \mathbb{C}$ is *not* an eigenvalue of D, and pick $u \in L_{\lambda}$, $u \neq 0$, $v \in R_{\lambda}$, $v \neq 0$. It is clear from this proof that $L_{\lambda} = \text{Span}\{u\}$ and $R_{\lambda} = \text{Span}\{v\}$. As $L_{\lambda} \cap R_{\lambda} = E_{\lambda} = \{0\}$, the vectors u and v must be linearly independent. Since they both lie in S_{λ} , which is 2-dimensional by Corollary 5.2.2, we can then conclude that $\{u, v\}$ is a basis for S_{λ} .

5.4 A special case

We go back to the Sturm-Liouville problem for the equation $D(y) = \lambda y$, $y \in Y$. As long as we are not able to show that D has eigenvalues, it is not possible for us to make efficient use of its spectral properties. Ideally, we would like to show that D is diagonalizable, in the sense that there exists a sequence of eigenfunctions of D in Y which forms an orthonormal basis for $L^2([a, b])$. The trick to make progress on this problem is to turn our attention to the inverse of D, whenever this makes sense.

We therefore assume throughout this section that $D: Y \to C([a, b])$ is 1-1. We will see how to get rid of this assumption in the next section.

We will first show that D(Y) = C([a, b]), i.e., D is onto, and that the inverse operator

$$D^{-1}: C([a,b]) \to Y$$

is an integral operator associated to a continuous kernel $G: [a, b] \times [a, b] \to \mathbb{C}$.

It should be noted here that, given a function $f \in C([a, b])$, the standard way to show that the differential equation -y'' + qy = f has a solution is to pick a basis for the associated homogeneous equation and use the method called *variation of parameters*. We will not discuss this method here and follow a shorter path.

Since the operator D is linear, the fact that D is 1-1 means that its kernel is trivial, that is, 0 is *not* an eigenvalue of D. As we saw in Remark 5.3.4, we can then pick a basis $\{u, v\}$ for

$$S_0 = \{y \in C^2([a,b]) : D(y) = 0\} = \{y \in C^2([a,b]) : y'' = qy\}$$

such that

- u satisfies the condition $\alpha_1 u(a) + \alpha_2 u'(a) = 0$,
- v satisfies the condition $\beta_1 v(b) + \beta_2 v'(b) = 0.$

Since q is real-valued, we can also assume that u and v are real-valued, cf. Remark 5.2.4.

We note that Remark 5.2.5 tells us that $W_{u,v}(x) \neq 0$ for all $x \in [a, b]$. Moreover, as $\widetilde{D}(u) = \widetilde{D}(v) = 0$, Lemma 5.2.6 gives that

$$(W_{u,v})' = \widetilde{D}(u)v - \widetilde{D}(v) = 0.$$

Hence, $W_{u,v}$ is a constant function on [a, b]. This means that

$$W_{u,v}(x) = u(x)v'(x) - v(x)u'(x) = W \quad \text{for all } x \in [a,b]$$

for some $W \in \mathbb{R} \setminus \{0\}$.

We can now define the associated Green's function $G: [a, b] \times [a, b] \to \mathbb{R}$ by

$$G(x,t) = -\frac{1}{W} \cdot \begin{cases} u(x)v(t) & \text{if } a \le x \le t \le b, \\ u(t)v(x) & \text{if } a \le t \le x \le b. \end{cases}$$

It is then straightforward to see that G is continuous. Hence we may form the associated integral operator $T_G: L^2([a,b]) \to L^2([a,b])$, which is given by

$$[T_G(f)](x) = \int_a^b G(x,t) f(t) dt$$

for all $f \in L^2([a,b])$ and $x \in [a,b]$. It is clear that T_G maps C([a,b]) into itself. In fact, it maps C([a,b]) into Y:

Proposition 5.4.1. Let $f \in C([a, b])$ and set $y := T_G(f)$. Then $y \in Y$. Moreover, D(y) = f.

Proof. Let $x \in [a, b]$. Using the definitions of G and T_G we get

$$y(x) = -\int_{a}^{x} W^{-1} v(x)u(t)f(t) dt - \int_{x}^{b} W^{-1} u(x)v(t)f(t) dt.$$

This implies that -Wy(x) = v(x)A(x) + u(x)B(x), where

$$A(x) := \int_a^x u(t)f(t) dt \quad \text{and} \quad B(x) := \int_x^b v(t)f(t) dt.$$

Thus we get

$$-Wy'(x) = v'(x)A(x) + v(x)A'(x) + u'(x)B(x) + u(x)B'(x)$$

= $v'(x)A(x) + v(x)u(x)f(x) + u'(x)B(x) - u(x)v(x)f(x)$
= $v'(x)A(x) + u'(x)B(x)$.

Since v', A, u' and B all lie in $C^1([a, b])$, we see that $y \in C^2([a, b])$.

Further, using that A(a) = 0 and $\alpha_1 u(a) + \alpha_2 u'(a) = 0$, we get

$$\alpha_1 y(a) + \alpha_2 y'(a) = \frac{-1}{W} \Big(\alpha_1 v(a) A(a) + \alpha_1 u(a) B(a) + \alpha_2 v'(a) A(a) + \alpha_2 u'(a) B(a) \Big)$$

= $\frac{-1}{W} \Big(\alpha_1 u(a) + \alpha_2 u'(a) \Big) B(a) = 0.$

In a similar way we get $\beta_1 v(b) + \beta_2 v'(b) = 0$. Thus we have shown that $y \in Y$.

To verify the second assertion, we first compute -Wy''. Since u'' = qu, v'' = qv, and v'u - u'v = W on [a, b], we get

$$-Wy'' = (v'A + u'B)' = v''A + v'A' + u''B + u'B'$$

= $q(vA + uB) + (v'u - u'v)f$
= $-qWy + (v'u - u'v)f$
= $W(f - qy).$

Thus, -y'' = f - qy, which gives

$$D(y) = -y'' + qy = f - qy + qy = f,$$

as desired.

The first part of Proposition 5.4.1 shows that D is onto C([a, b]). Since D is also 1-1 (by assumption), D has an inverse map $D^{-1}: C([a, b]) \to Y$, which is defined as follows:

Given some $f \in C([a, b])$, then

$$D^{-1}(f) := y \,,$$

where $y \in Y$ is the unique function in Y such that D(y) = f. We now see that the second part of Proposition 5.4.1 tells us that

$$D^{-1}(f) = T_G(f)$$
 for every $f \in C([a, b])$, i.e., $D^{-1} = (T_G)_{|C([a, b])}$.

Since G(x,t) = G(t,x) for all (x,t) in $[a,b] \times [a,b]$ (check this!), we get that T_G is self-adjoint. As T_G is also compact (indeed, it is a Hilbert-Schmidt operator on $L^2([a,b])$, cf. Example 4.2.11), we are in the position to apply the spectral theorem to T_G .

However, we will also need to know that T_G maps $L^2([a, b])$ into C([a, b]). This is true for any integral operator with continuous kernel:

Lemma 5.4.2. Assume $K : [a,b] \times [a,b] \to \mathbb{C}$ is continuous. Then the associated integral operator $T_K : L^2([a,b]) \to L^2([a,b])$ maps $L^2([a,b])$ into C([a,b]).

Proof. Let $f \in L^2([a, b])$ and let $\varepsilon > 0$. Note that the Cauchy-Schwarz inequality gives that

$$M := \int_{[a,b]} |f| \, dm \le \left(\int_{[a,b]} 1 \, dm \right)^{1/2} \left(\int_{[a,b]} |f|^2 \, dm \right)^{1/2} = \sqrt{b-a} \, \|f\|_2 < \infty,$$

where *m* denotes the Lebesgue measure on [a, b]. As the continuous function *K* is automatically uniformly continuous on the compact set $R := [a, b] \times [a, b]$, we can find $\delta > 0$ such that

$$\left|K(x_1,t_1) - K(x_2,t_2)\right| < \varepsilon/M$$

whenever $(x_1, t_1), (x_2, t_2) \in R$ and $|x_2 - x_1| < \delta, |t_2 - t_1| < \delta.$

Let now $x_0 \in [a, b]$. Then for every $t \in [a, b]$ and all $x \in [a, b]$ such that $|x - x_0| < \delta$, we have

$$|K(x,t) - K(x_0,t)| < \varepsilon/M.$$

Thus we get

$$\left| [T_K(f)](x) - [T_K(f)](x_0) \right| = \left| \int_{[a,b]} \left(K(x,t) - K(x_0,t) \right) f(t) \, dm(t) \right|$$
$$\leq \int_{[a,b]} \left| K(x,t) - K(x_0,t) \right| |f(t)| \, dm(t)$$
$$\leq \varepsilon / M \int_{[a,b]} |f| \, dm = \varepsilon$$

for all $x \in [a, b]$ such that $|x - x_0| < \delta$. This shows that $T_K(f)$ is continuous at x_0 .

Since x_0 was an arbitrary point of $[a, b], T_K(f) \in C([a, b])$.

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Theorem 5.4.3. Assume that D is 1-1, and consider the Sturm-Liouville problem

$$D(y) = \lambda y$$
 with $y \in Y$.

Then the following assertions hold:

- The eigenvalues for this problem form a countable set $\{\lambda_k : k \in \mathbb{N}\}$ of non-zero distinct real numbers satisfying that $|\lambda_k| \to \infty$ as $k \to \infty$.
- For each $k \in \mathbb{N}$ the eigenspace $E_{\lambda_k} = \{y \in Y : D(y) = \lambda_k y\}$ is one-dimensional.
- If y_k is a unit vector in E_{λ_k} for each $k \in \mathbb{N}$, then $\{y_k : k \in \mathbb{N}\}$ is an orthonormal basis for $L^2([a, b])$.

Proof. We first observe that 0 is not an eigenvalue of T_G :

Indeed, since Y is dense in $H := L^2([a, b])$, and $Y = T_G(C([a, b]))$, we have

$$H = \overline{Y} \subset \overline{T_G(H)} \subset H,$$

hence $\overline{T_G(H)} = H$. Thus, we get $\ker(T_G)^{\perp} = \overline{T_G(H)} = H$, i.e., $\ker(T_G) = \{0\}$.

Applying the spectral theorem to T_G , we obtain that the eigenvalues of T_G form a countable set $\{\mu_k : k \in \mathbb{N}\}$ of non-zero distinct real numbers satisfying that $\mu_k \to 0$ as $k \to \infty$.

Let $k \in \mathbb{N}$, and set $\lambda_k = \mu_k^{-1} \neq 0$. Let $f_k \in H$ be an eigenfunction for T_G associated to μ_k . Since $T_G(f_k) = \mu_k f_k$, we get

$$f_k = \lambda_k T_G(f_k) \,. \tag{5.4.1}$$

As Lemma 5.4.2 gives that $T_G(f_k) \in C([a, b])$, this gives that $f_k \in C([a, b])$. Hence, $T_G(T_G(f_k)) \in Y$. But (5.4.1) implies that

$$f_k = \lambda_k^2 T_G(T_G(f_k)),$$

so we get that $f_k \in Y$. Now, applying D to (5.4.1), we get

$$D(f_k) = \lambda_k f_k$$

This shows that λ_k is an eigenvalue of D, and f_k is an eigenfunction for D associated to λ_k . Now, Proposition 5.3.3 tells us that $E_{\lambda_k} := \{y \in Y : D(y) = \lambda_k y\}$ is one-dimensional. Hence, we have $E_{\lambda_k} = \text{Span}\{f_k\}$. Further, one readily checks that E_{λ_k} is also the eigenspace of T_G associated to μ_k .

We note that D can not have other eigenvalues than the λ_k 's (for if D had one such eigenvalue, then T_G would have an eigenvalue different from all μ_k 's, which is not the case). Further, we note that

$$\lim_{k \to \infty} |\lambda_k| = \lim_{k \to \infty} |\mu_k|^{-1} = \infty$$

since $\lim_{k\to\infty} \mu_k = 0$.

Finally, if we set $y_k := \pm (||f_k||_2)^{-1} f_k \in E_{\lambda_k}$ for each $k \in \mathbb{N}$, then we also get from the spectral theorem for T_G that $\{y_k : k \in \mathbb{N}\}$ is an orthonormal basis for H consisting of eigenfunctions for D.

Example 5.4.4. To illustrate this theorem, let us consider the Sturm-Liouville system

$$-y'' = \lambda y$$
 on $[0, \pi]$, with $Y = \{y \in C^2([0, \pi]) : y(0) = y(\pi) = 0\}.$

Thus we have $D(y) = -y'', y \in Y$.

Let $\lambda \in \mathbb{C}$. We consider first the case $\lambda = 0$. It is straightforward to check that the only function y in Y satisfying -y'' = 0 is y = 0. Thus, 0 is not an eigenvalue of D, i.e., D is 1-1, so Theorem 5.4.3 applies in this case. We can determine the eigenvalues and the eigenfunctions of D explicitly as follows.

Assume $\lambda \neq 0$, and write $\lambda^{1/2} = r + i s$ with $(r, s) \in \mathbb{R}^2 \setminus \{(0, 0)\}$.³ A basis for $S_{\lambda} := \{y \in C^2([0, \pi]) : -y'' = \lambda y\}$ is given by

$$\{e^{i\lambda^{1/2}x}, e^{-i\lambda^{1/2}x}\} = \{e^{-sx}(\cos(rx) + i\sin(rx)), e^{sx}(\cos(rx) - i\sin(rx))\}.$$

If $y \in S_{\lambda}$, say $y(x) = C_1 e^{-sx}(\cos(rx) + i(sx)) + C_2 e^{sx}(\cos(rx) - i\sin(rx))$, then $y \in Y$ if and only if

$$\begin{cases} C_1 + C_2 = 0, \\ C_1 e^{-s\pi} (\cos(r\pi) + i\,\sin(r\pi)) + C_2 e^{s\pi} (\cos(r\pi) - i\,\sin(r\pi)) = 0. \end{cases}$$

This gives that $E_{\lambda} = S_{\lambda} \cap Y$ is non-trivial if and only if

$$e^{-s\pi}(\cos(r\pi) + i\,\sin(r\pi)) = e^{s\pi}(\cos(r\pi) - i\,\sin(r\pi))$$

and it is elementary to deduce that this happens if and only if s = 0 and r = k for some $k \in \mathbb{Z} \setminus \{0\}$, in which case $\lambda = k^2 \in \mathbb{N}$ and $E_{\lambda} = \text{Span}\{\sin(kx)\}$.

³We could here have used that we know that all the possible eigenvalues of D are real, so that we need only to consider $\lambda \in \mathbb{R} \setminus \{0\}$. However this would not shorten our discussion significantly.

5. An introduction to Sturm-Liouville theory

This means that the distinct eigenvalues of this Sturm-Liouville system are $\lambda_k = k^2, k \in \mathbb{N}$, with associated normalized eigenfunctions $y_k(x) = \sqrt{2/\pi} \sin(kx)$. This is in accordance with Theorem 5.4.3. Note that this theorem implies that the set $\{\sqrt{2/\pi} \sin(kx) : k \in \mathbb{N}\}$ is an orthonormal basis for $L^2([0,\pi])$.

For completeness, we also compute the Green's function G and $T_G(f)$ for $f \in C([0, \pi])$. One computes easily that

$$L_0 = \{ y \in C^2([0,\pi]) : y'' = 0, \ y(0) = 0 \} = \operatorname{Span}\{u\},\$$

where u(x) = x, while

$$R_0 = \{ y \in C^2([0,\pi]) : y'' = 0, \ y(\pi) = 0 \} = \operatorname{Span}\{v\},\$$

where $v(x) = x - \pi$. Thus we get that

$$W = W_{u,v}(x) = u(x)v'(x) - v(x)u'(x) = x - (x - \pi) = \pi$$

for all $x \in [0, \pi]$. Moreover, the Green's function $G : [0, \pi] \times [0, \pi] \to \mathbb{C}$ is then given by

$$G(x,t) = \frac{1}{\pi} \cdot \begin{cases} x \left(\pi - t \right) & \text{if } 0 \le x \le t \le \pi, \\ t \left(\pi - x \right) & \text{if } 0 \le t \le x \le \pi, \end{cases}$$

and we obtain that

$$[T_G(f)](x) = \int_0^{\pi} G(x,t) f(t) dt = \frac{1}{\pi} \left((\pi - x) \int_0^x tf(t) dt + x \int_x^{\pi} (\pi - t)f(t) dt \right)$$

for all $f \in C([0,\pi])$ and $x \in [0,\pi]$.

Note that determining the eigenvalues of T_G by direct computation is not an easy task. Anyhow, you should verify that $\sin(kx)$ is an eigenfunction for T_G associated with the eigenvalue $\mu_k = k^{-2}$ for each $k \in \mathbb{N}$.

5.5 The general case

In this final section, we consider the general case, i.e., we don't assume that D is 1-1. The idea now is to show that there exists some $\mu \in \mathbb{R}$ which is not an eigenvalue of D, and consider the operator $D_{\mu} : Y \to C([a, b])$ defined by $D_{\mu}(y) = -y'' + qy - \mu y$. Then 0 will not be an eigenvalue of D_{μ} (otherwise there would be some $y \in Y \setminus \{0\}$ such that $D_{\mu}(y) = 0$, i.e.,

 $D(y) = \mu y$, and μ would be an eigenvalue of D, giving a contradiction). Hence, we will be able to apply Theorem 5.4.3 to the Sturm-Liouville system $D_{\mu}(y) = \lambda' y$ with $y \in Y$, and deduce some interesting consequences for our original Sturm-Liouville problem.

Since we know that the possible eigenvalues of D are all real numbers, we may think: why not just pick some $\mu \in \mathbb{C} \setminus \mathbb{R}$? The problem with such a choice is that the function $q_{\mu}(x) := q(x) - \mu$ for $x \in [a, b]$ will not be real-valued, hence that the Sturm-Liouville system associated with D_{μ} will not match our requirements.

Lemma 5.5.1. There exists some $\mu \in \mathbb{R}$ which is not an eigenvalue of D.

Proof. We know that $L^2([a, b])$ has a countable orthonormal basis, say $\{u_k\}_{k\in\mathbb{N}}$.⁴ We will show that this implies that D has a countable number of distinct eigenvalues. Since \mathbb{R} is uncountable, the assertion to be proven will clearly follow.

Assume (for contradiction) that D has an uncountable number of distinct eigenvalues. Then we can pick a unit vector in each of the associated eigenspaces. As all eigenspaces of D are orthogonal to each other, this means that there exists an orthonormal subset Γ of Y which is uncountable.

Let $k \in \mathbb{N}$. Then it follows from Bessel's inequality that

$$M_k := \sup_{A \subset \Gamma, A \text{ finite}} \left\{ \sum_{\gamma \in A} |\langle u_k, \gamma \rangle|^2 \right\} \le ||u_k||_2^2 = 1 < \infty.$$

This implies that the set

$$U_{k,n} := \left\{ \gamma \in \Gamma : |\langle u_k, \gamma \rangle| \ge \frac{1}{n} \right\}$$

is finite for every $n \in \mathbb{N}$: indeed, if $U_{k,n}$ was infinite for some n, then we could find an infinite sequence $\{\gamma_m\}_{m\in\mathbb{N}}$ of distinct elements in $U_{k,n}$, and this would give that

$$M_k \geq \sup_{m \in \mathbb{N}} \left\{ \sum_{\ell=1}^m |\langle u_k, \gamma_\ell \rangle|^2 \right\} \geq \sup_{m \in \mathbb{N}} \left\{ m \cdot \frac{1}{n^2} \right\} = \infty,$$

contradicting that $M_k < \infty$. Setting now $U_k := \{\gamma \in \Gamma : \langle u_k, \gamma \rangle \neq 0\}$, we get that

$$U_k = \bigcup_{n \in \mathbb{N}} U_{k,n}$$

⁴One may for example take $u_k(x) = \sqrt{2/(b-a)} \sin \left(k\pi (x-a)/(b-a) \right)$ for $x \in [a,b]$ and $k \in \mathbb{N}$, cf. Example 5.4.4 when $[a,b] = [0,\pi]$.

is countable (being a countable union of finite sets). Hence, the countable union $U := \bigcup_{k \in \mathbb{N}} U_k$ is a countable subset of Γ . As Γ is uncountable, there must exist some $\gamma \in \Gamma \setminus U$. But then $\gamma \notin U_k$ for every $k \in \mathbb{N}$, so we have

$$\langle u_k, \gamma \rangle = 0$$
 for every $k \in \mathbb{N}$.

This says that γ is orthogonal to every u_k , so we must have $\gamma = 0$ (since $\{u_k\}_{k \in \mathbb{N}}$ is an orthonormal basis for $L^2([a, b])$). But this gives a contradiction, since every element of Γ is a unit vector. We can therefore conclude that D has a countable number of distinct eigenvalues.

We can now state our main result about regular Sturm-Liouville systems: **Theorem 5.5.2.** Consider the Sturm-Liouville problem

 $D(y) = \lambda y$ with $y \in Y$.

Then the following assertions hold:

- The eigenvalues for this problem form a countable set {λ_k : k ∈ N} of distinct real numbers satisfying that |λ_k| → ∞ as k → ∞.
- For each $k \in \mathbb{N}$ the eigenspace $E_{\lambda_k} = \{y \in Y : D(y) = \lambda_k y\}$ is one-dimensional.
- If y_k is a unit vector in E_{λ_k} for each $k \in \mathbb{N}$, then $\{y_k : k \in \mathbb{N}\}$ is an orthonormal basis for $L^2([a, b])$.

Proof. By Lemma 5.5.1 we can find some $\mu \in \mathbb{R}$ which is not an eigenvalue of D. We define $D_{\mu}: Y \to C([a, b])$ by $D_{\mu}(y) = -y'' + qy - \mu y$ and consider the Sturm-Liouville system $D_{\mu}(y) = \lambda' y$ on Y. Then, as 0 is not an eigenvalue of D_{μ} , we may apply Theorem 5.4.3 to D_{μ} . This gives:

- The eigenvalues of D_{μ} form a countable set $\{\lambda'_k : k \in \mathbb{N}\}$ of non-zero distinct real numbers satisfying that $|\lambda'_k| \to \infty$ as $k \to \infty$.
- For each $k \in \mathbb{N}$ the eigenspace $E'_{\lambda'_k} = \{y \in Y : D_{\mu}(y) = \lambda'_k y\}$ is one-dimensional.
- If v_k is a unit vector in $E'_{\lambda'_k}$ for each $k \in \mathbb{N}$, then $\{v_k : k \in \mathbb{N}\}$ is an orthonormal basis for $L^2([a, b])$.

Now, for $y \in Y$, we obviously have $D(y) = \lambda y$ if and only if $D_{\mu}(y) = (\lambda - \mu)y$. This implies that the set consisting of all eigenvalues of D is the countable set of distinct real numbers given by $\{\lambda_k : k \in \mathbb{N}\}$, where $\lambda_k := \lambda'_k + \mu$ for each $k \in \mathbb{N}$. Moreover, the eigenspace E_{λ_k} of D associated to each λ_k is then equal to $E'_{\lambda'_k}$, hence is one-dimensional.

Finally, if y_k is a unit vector in E_{λ_k} for each $k \in \mathbb{N}$, then we have $y_k = \pm v_k$ for every $k \in \mathbb{N}$, so the last assertion clearly follows.

5.6 Exercises

Exercise 5.1. Find the eigenvalues and eigenfunctions of the Sturm-Liouville system $-y'' = \lambda y$ on the given interval with the following boundary conditions:

- a) $[a,b] = [0,\pi], y'(0) = 0, y'(\pi) = 0.$
- b) $[a,b] = [0,\pi], y'(0) = 0, y(\pi) = 0.$
- c) $[a,b] = [0,2\pi], y(0) = 0, y(2\pi) = 0.$
- d) [a,b] = [0,1], y(0) = 0, y(1) + y'(1) = 0.

Exercise 5.2. Consider a Sturm-Liouville system $D(y) = \lambda y$ on Y as in (5.1.4), with boundary conditions as in (5.1.2).

Assume that the following two extra conditions holds:

- a) $q(x) \ge 0$ for all $x \in [a, b]$,
- b) $\alpha_1 \alpha_2 \leq 0$ and $\beta_1 \beta_2 \geq 0$.

Show that the eigenvalues of D are all non-negative.

Exercise 5.3. Consider a Sturm-Liouville system $D(y) = \lambda y$ as in (5.1.4), but where

$$Y = \{ y \in C^2([a, b]) \mid y(a) \text{ or } y'(a) = 0 ; y(b) \text{ or } y'(b) = 0 \}.$$

Show that the distinct eigenvalues of D may ordered so that

$$\lambda_1 < \lambda_2 < \lambda_3 < \cdots$$
 and $\lim_{k \to \infty} \lambda_k = \infty$.

Exercise 5.4. Assume D is 1-1, as in section 4. Show that the distinct eigenvalues of D satisfy

$$\sum_{k=1}^{\infty} \, rac{1}{|\lambda_k|^2} \, < \, \infty \, .$$

Exercise 5.5. Show that the space

 $Y = \{y \in C^2([a, b]) : y \text{ satisfies the conditions (5.1.2)}\}$ is dense in $L^2([a, b])$.